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## 0 Measure Theory

#### 0.1 Key Theorems

#### Theorem 0.1. *Fubini's Theorem* http://en.wikipedia.org/wiki/Fubini%27s\_theorem

Suppose A and B are complete measure spaces. Suppose f(x, y) is  $A \times B$  measurable. If

$$\int_{A\times B} |f(x,y)| \, d(x,y) < \infty$$

where the integral is taken with respect to a product measure on the space over  $A \times B$ , then

$$\int_{A} \left( \int_{B} f(x, y) \, dy \right) \, dx = \int_{B} \left( \int_{A} f(x, y) \, dx \right) \, dy = \int_{A \times B} f(x, y) \, d(x, y)$$

the first two integrals being iterated integrals with respect to two measures, respectively, and the third being an integral with respect to a product of these two measures.

#### **Corollary:**

If f(x,y) = g(x)h(y) for some functions g and h, then

$$\int_{A} g(x) \, dx \int_{B} h(y) \, dy = \int_{A \times B} f(x, y) \, d(x, y)$$

the third integral being with respect to a product measure.

#### Theorem 0.2. *Tonelli's Theorem* http://en.wikipedia.org/wiki/Fubini%27s\_theorem#Tonelli.27s\_theorem

Suppose that A and B are  $\sigma$ -finite measure spaces, not necessarily complete. If either

$$\int_{A} \left( \int_{B} |f(x,y)| \, dy \right) \, dx < \infty \text{ or } \int_{B} \left( \int_{A} |f(x,y)| \, dx \right) \, dy < \infty$$

then

$$\int_{A \times B} |f(x,y)| \, d(x,y) < \infty$$

and

$$\int_{A} \left( \int_{B} f(x,y) \, dy \right) \, dx = \int_{B} \left( \int_{A} f(x,y) \, dx \right) \, dy = \int_{A \times B} f(x,y) \, d(x,y)$$

#### Remark 0.3. Fubini vs. Tonelli

http://en.wikipedia.org/wiki/Fubini%27s\_theorem

Tonelli's theorem is a successor of Fubini's theorem. The conclusion of Tonelli's theorem is identical to that of Fubini's theorem, but the assumptions are different. Tonelli's theorem states that on the product of two -finite measure spaces, a product measure integral can be evaluated by way of an iterated integral for nonnegative measurable functions, regardless of whether they have finite integral. A formal statement of Tonelli's theorem is identical to that of Fubini's theorem, except that the requirements are now that  $(X, A, \mu)$  and  $(Y, B, \nu)$  are  $\sigma$ -finite measure spaces, while f maps  $X \times Y$  to  $[0, \infty]$ .

Theorem 0.4. *Hölder's Inequality* Theorem 12.54 on page 356

Let  $1 \le p, q \le \infty$  and  $\frac{1}{p} + \frac{1}{q} = 1$ . If  $f \in L^p(X, \mu)$  and  $g \in L^q(X, \mu)$ , then  $fg \in L^1(X, \mu)$  and

$$\int fg \, d\mu \bigg| \le \|f\|_p \|g\|_q$$

Theorem 0.5. Young's Inequality Theorem 12.58 on page 359

Let  $1 \leq p, q, r \leq \infty$  and  $\frac{1}{p} + \frac{1}{q} = 1 + \frac{1}{r}$ . If  $f \in L^p(\mathbb{R}^n)$  and  $g \in L^q(\mathbb{R}^n)$ , then  $f * g \in L^r(\mathbb{R}^n)$  and

 $||f * g||_r \le ||f||_p ||g||_q$ 

Theorem 0.6. *Lebesgue Dominated Convergence Theorem* Theorem 12.35 on page 348

Suppose that  $(f_n)$  is a sequence of integrable functions,  $f_n : X \to \overline{\mathbb{R}}$ , on a measure space  $(X, \mathcal{A}, \mu)$  that converges pointwise to a limiting function  $f : X \to \overline{\mathbb{R}}$ . If there is an integrable function  $g : X \to [0, \infty]$  such that

 $|f_n(x)| \le g(x) \quad \forall \ x \in X, \ n \in \mathbb{N}$ 

then f is integrable and

$$\lim_{n \to \infty} \int f_n \, d\mu = \int f \, d\mu$$

Theorem 0.7. Cauchy-Schwarz Inequality http://en.wikipedia.org/wiki/Cauchy-Schwarz\_inequality

**Formal Statement:** For all vectors x, y of an inner product space,

$$\begin{aligned} \left| \langle x, y \rangle \right|^2 &\leq \langle x, x \rangle \left\langle y, y \right\rangle \\ \left| \langle x, y \rangle \right| &\leq \|x\| \|y\| \end{aligned}$$

Square of a Sum:

$$\left|\sum_{i=1}^{n} x_i y_i\right|^2 \le \sum_{i=1}^{n} |x_i|^2 \sum_{i=1}^{n} |y_i|^2$$

In  $L^2$ :

$$\left|\int f(x)g(x)\,dx\right|^2 \le \int |f(x)|^2\,dx \int |g(x)|^2\,dx$$

# 7 Fourier Series

#### 7.1 Fourier Series

**Definition 7.1.**  $2\pi$ -*periodic* page 149

A function  $f : \mathbb{R} \to \mathbb{C}$  is  $2\pi$ -periodic if

$$f(x+2\pi) = f(x) \quad \forall \ x \in \mathbb{R}$$

A  $2\pi$ -periodic function may be indentified with a function on the unit circle, or one-dimensional torus,  $\mathbb{T} = \mathbb{R}/(2\pi\mathbb{Z})$ . The space  $C(\mathbb{T})$  is the space of continuous functions from  $\mathbb{T}$  to  $\mathbb{C}$ , and  $L^2(\mathbb{T})$  is the completion of  $C(\mathbb{T})$  with respect to the  $L^2$ -norm,

$$||f|| = \left(\int_{\mathbb{T}} |f(x)|^2 \, dx\right)^{1/2}$$

 $L^2(\mathbb{T})$  is a Hilbert space with respect to the inner product

$$\langle f,g \rangle = \int_{\mathbb{T}} \overline{f(x)} g(x) \, dx$$

**Definition 7.2.**  $L^p(\mathbb{T})$ page 92 and Notes 1/3/11

 $L^p(\mathbb{T}) :=$  the space of Lebesgue measurable functions,  $f : \mathbb{T} \to \mathbb{C}$  such that  $\int_{\mathbb{T}} |f|^p dx < \infty$ , where  $1 \le p < \infty$ . We define the  $L^p$ -norm as:

$$\|f\|_p = \left(\int_{\mathbb{T}} |f|^p \, dx\right)^{1/p}$$

For  $p = \infty$ ,  $L^{\infty}(\mathbb{T})$  is the space of Lebesgue measurable functions that are essentially bounded on  $\mathbb{T}$ , meaning that f is bounded on every subset of  $\mathbb{T}$  with nonzero measure. The norm on  $L^{\infty}(\mathbb{T})$  is the essential supremum

$$||f||_{\infty} = \inf\{M \mid |f(x)| \le M \text{ a.e. in } \mathbb{T}\}\$$

We identify f with g if f = g a.e. (almost everywhere, except possibly on a set of measure 0).

# Theorem 7.3.

Notes 1/3/11

 $L^p(\mathbb{T})$  with the norm  $\|f\|_{L^p} = (\int_{\mathbb{T}} |f|^p \, dx)^{1/p}$  is a Banach space.

 $C(\mathbb{T})$  is dense in  $L^p(\mathbb{T})$  for  $1 \leq p < \infty$ .

Note:  $C(\mathbb{T}) :=$  the space of continuous functions  $f : \mathbb{T} \to \mathbb{C}$ 

**Proposition 7.5.** Notes 1/3/11

 $p > q \Rightarrow L^p(\mathbb{T}) \subset L^q(\mathbb{T}) \text{ and } \|\cdot\|_p \ge \|\cdot\|_q$ Also,  $L^1(\mathbb{T}) \supset L^2(\mathbb{T}) \supset \ldots \supset C(\mathbb{T})$ 

**Example 7.6.** Fourier Basis Example Notes 1/3/11

$$\sum_{n \neq 0} \frac{1}{|n|} e^{inx} = f(x)$$
$$\sum_{n \neq 0} \frac{1}{|n|^2} = 2 \sum_{n=1}^{\infty} \frac{1}{n^2} < \infty$$
$$\lim_{N \to \infty} \int \left| f(x) - \sum_{n=-N, n \neq 0}^{N} \frac{1}{|n|} e^{inx} \right|^2 dx = 0$$

Line 2 and Bessel's Inequality tell us that the series converges in  $L^2(\mathbb{T})$ . However, it doesn't converge pointwise everywhere on  $\mathbb{T}$ .

Ex. at x = 0,  $\sum_{n \neq 0} \frac{1}{|n|}$  diverges.

**Proposition 7.7.** Orthonormal Basis of  $L^2(\mathbb{T})$ page 150

The Fourier basis elements are the functions

$$e_n(x) = \frac{1}{\sqrt{2\pi}} e^{inx}$$

 $\{e_n \mid n \in \mathbb{Z}\}\$  is an orthonormal basis of  $L^2(\mathbb{T})$ .

Proof Outline

• Orthogonality

It is easily shown that

$$\langle e_m, e_n \rangle = \begin{cases} 1 & \text{if } m = n \\ 0 & \text{if } m \neq n \end{cases}$$

• Completeness

This proof relies upon the ideas of convolution and approximate identities. (See Theorems 7.12 and 7.13.)

# **Definition 7.8.** Convolution

page 150

The *convolution* of two continuous functions  $f, g: \mathbb{T} \to \mathbb{C}$  is the continuous function  $f * g: \mathbb{T} \to \mathbb{C}$ defined by the integral

$$(f*g)(x) = \int_{\mathbb{T}} f(x-y)g(y) \, dy$$

Using the change of variable z = x - y, it is seen that

$$(f * g)(x) = \int_{\mathbb{T}} f(z)g(x - z) \, dz$$

so that f \* g = g \* f.

**Definition 7.9.** Approximate Identity Definition 7.1 on page 151

A family of functions  $\{\varphi_n \in C(\mathbb{T}) \mid n \in \mathbb{N}\}$  is an *approximate identity* if

(a) 
$$\varphi_n(x) \ge 0$$
  
(b)  $\int_{\mathbb{T}} \varphi_n(x) \, dx = 1$   
(c)  $\lim_{n \to \infty} \int_{\delta \le |x| \le \pi} \varphi_n(x) \, dx = 0 \quad \forall \ \delta > 0$ 

Note: in (c),  $\mathbb{T}$  is identified with  $[-\pi, \pi]$ .

#### Theorem 7.10.

Theorem 7.2 on page 151 and Notes 1/5/11 and FA 49

If  $\{\varphi_n \in C(\mathbb{T}) \mid n \in \mathbb{N}\}\$  is an approximate identity and  $f \in C(\mathbb{T})$ , then  $\varphi_n * f$  converges uniformly to f as  $n \to \infty$ .

Note: the term "approximate identity" comes from this result, since  $\{\varphi_n\}$  is an approximation to the identity.

Proof

$$f(x) = \int_{\mathbb{T}} \varphi_n(y) f(x) \, dy$$
$$(\varphi_n * f)(x) = \int_{\mathbb{T}} \varphi_n(y) f(x - y) \, dy$$
$$(\varphi_n * f)(x) - f(x) = \int_{\mathbb{T}} \varphi_n(y) [f(x - y) - f(x)] \, dy$$

- f is uniformly continuous, so there exists M such that  $|f(x)| \leq M \ \forall \ x \in \mathbb{T}$
- $\exists \ \delta > 0$  such that  $|f(x) f(y)| \le \epsilon$  whenever  $|x y| < \delta$

$$\begin{split} |(\varphi_n * f)(x) - f(x)| &\leq \int_{-\pi}^{\pi} \varphi_n(y) |f(x - y) - f(x)| \, dy \\ &\leq \int_{|y| < \delta} \varphi_n(y) |f(x - y) - f(x)| \, dy + \int_{|y| \ge \delta} \varphi_n(y) |f(x - y) - f(x)| \, dy \\ &\leq \epsilon \int_{|y| < \delta} \varphi_n(y) \, dy + \int_{|y| \ge \delta} \varphi_n(y) [|f(x - y)| + |f(x)|] \, dy \\ &\leq \epsilon + 2M \int_{|y| \ge \delta} \varphi_n(y) \, dy \end{split}$$

Using property (c) of an approximate identity gives that  $\varphi_n * f \to f$  uniformly in  $C(\mathbb{T})$ .

#### **Remark 7.11.** Revised Approximate Identity Definition Notes 1/5/11

More generally,  $\varphi_n \in L^1(\mathbb{T})$  is an approximate identity if

(a) 
$$\int_{\mathbb{T}} |\varphi_n(x)| \, dx \le M \quad \forall \ n \in \mathbb{N}$$
  
(b) 
$$\int_{\mathbb{T}} \varphi_n(x) \, dx = 1$$
  
(c) 
$$\lim_{n \to \infty} \int_{\delta \le |x| \le \pi} \varphi_n(x) \, dx = 0 \quad \forall \ \delta > 0$$

**Theorem 7.12.** Weierstrass Approximation Theorem Theorem 7.3 on page 152 and Notes 1/5/11

The trigonometric polynomials are dense in  $C(\mathbb{T})$  with respect to the uniform norm.

Proof

- Let  $f \in C(\mathbb{T})$
- Generate an approximate identity that is a trigonometric polynomial

- Define  $\varphi_n = c_n (1 + \cos x)^n = c_n [2\cos^2(\frac{x}{2})]^n$  and choose  $c_n$  such that  $\int_{\mathbb{T}} \varphi_n(x) \, dx = 1$
- To show  $\varphi_n$  is an approximate identity, we need to show that  $\forall \ \delta > 0$ ,  $\lim_{n \to \infty} \int_{|x| > \delta} \varphi_n(x) \, dx = 0$ 
  - \* Fix  $\epsilon > 0$ .  $\forall x, \delta \le |x| \le \pi, \exists r \in (0, 1)$  such that

$$(1 + \cos x) < r(1 + \cos y)$$
  

$$\varphi_n(x) < r^n \varphi_n(y)$$
  

$$\delta \varphi_n(x) < r^n \int_{-\delta/2}^{\delta/2} \varphi_n(y) \, dy$$
  

$$\delta \varphi_n(x) < r^n$$

$$0 \le \varphi_n(x) < \frac{r^n}{\delta} \quad \forall x \text{ such that } \delta \le |x| \le \pi$$

- So  $\varphi_n \to 0$  uniformly on  $\delta \le |x| \le \pi$  as  $n \to \infty$ , and  $\int_{|x| > \delta} \varphi_n(x) \, dx \to 0$  as  $n \to \infty$
- $\varphi_n$  is an approximate identity, so  $\varphi_n * f$  is a trigonometric polynomial, and  $\varphi_n * f$  converges uniformly to f (See Theorem 7.10)

## Corollary 7.13.

page 153 and 155 and Notes 1/5/11

The trigonometric polynomials are dense in  $L^2(\mathbb{T})$ . That is, for any  $f \in L^2(\mathbb{T})$ ,

$$f(x) = \frac{1}{\sqrt{2\pi}} \sum_{n=-\infty}^{\infty} \hat{f}_n e^{inx}$$

$$\hat{f}_n = \frac{1}{\sqrt{2\pi}} \int_{\mathbb{T}} f(x) e^{-inx} \, dx$$

If  $f \in L^2(\mathbb{T})$  then the Fourier series of f converges pointwise to f a.e. (Carleson).

 $\frac{\text{Proof}}{\text{Let } f} \in L^2(\mathbb{T}).$ 

- Choose  $g \in C(\mathbb{T})$  such that  $||f g||_{L^2} < \epsilon/2$ . We can do this because  $C(\mathbb{T})$  is dense in  $L^2(\mathbb{T})$ .
- Pick a trigonometric polynomial p such that  $||g p||_{L^2} < \epsilon/2\sqrt{2\pi}$ .

• 
$$||g - p||_{L^2} = (\int |g - p|^2 dx)^{1/2} \le ||g - p||_{\infty} \sqrt{2\pi}$$

•  $||f - p||_{L^2} \le ||f - g||_{L^2} + ||g - p||_{L^2} < \epsilon/2 + \epsilon/2$ 

Corollary 7.14. Notes 1/5/11

 $\{\frac{1}{\sqrt{2\pi}}e^{inx} \mid n \in \mathbb{Z}\}$  is an orthonormal basis of  $L^2(\mathbb{T})$ .

# **Definition 7.15.** *Periodic Fourier Transform* page 153 and Notes 1/7/11

The *Periodic Fourier Transform*  $\mathcal{F} : L^2(\mathbb{T}) \to \ell^2(\mathbb{Z})$  maps a function to its sequence of Fourier coefficients by

$$\mathcal{F}f = \left(\hat{f}_n\right)_{n=-\infty}^{\infty}$$

Thus, the  $L^2$  norm of a function can be computed by

$$\int_{\mathbb{T}} |f(x)|^2 \, dx = \sum_{n=-\infty}^{\infty} \left| \hat{f}_n \right|^2$$

This implies that  $(\hat{f}_n) \in \ell^2(\mathbb{Z})$ . Furthermore, the Projection Theorem (6.13 in the book) implies that

$$f_N(x) = \frac{1}{\sqrt{2\pi}} \sum_{n=-N}^{N} \hat{f}_n e^{inx}$$

is the best approximation of f by a trigonometric polynomial of degree N in the  $L^2$ -norm.

Theorem 7.16. *Parseval's Theorem* Notes 1/7/11

Given  $f, g \in L^2(\mathbb{T})$ , then

$$f(x) = \frac{1}{\sqrt{2\pi}} \sum_{n \in \mathbb{Z}} a_n e^{inx}$$
$$g(x) = \frac{1}{\sqrt{2\pi}} \sum_{n \in \mathbb{Z}} b_n e^{inx}$$
$$\langle f, g \rangle = \sum_{n \in \mathbb{Z}} \overline{a_n} b_n$$

#### Proposition 7.17.

Proposition 7.4 on page 154

If 
$$f, g \in L^2(\mathbb{T})$$
, then  $f * g \in C(\mathbb{T})$  and

 $||f * g||_{\infty} \le ||f||_2 ||g||_2$ 

Proof

$$(f * g)(x) = \int_{\mathbb{T}} f(x - y)g(y) \, dy$$

If  $f, g \in C(\mathbb{T})$ , then we can apply the Cauchy-Schwarz Inequality to get

$$|f * g(x)| \le ||f||_{L^2} ||g||_{L^2}$$

Taking the supremum of both sides yields

$$\|f * g\|_{\infty} \le \|f\|_{L^2} \|g\|_{L^2}$$

If  $f, g \in L^2(\mathbb{T})$ , then there exist sequences  $(f_k), (g_k) \in C(\mathbb{T})$  such that  $||f - f_k||_2 \to 0$  and  $||g - g_k||_2 \to 0$  as  $k \to \infty$ . Also, the sequence  $(f_k * g_k) \in C(\mathbb{T})$  is Cauchy with respect to the sup-norm, since

$$\begin{aligned} \|f_j * g_j - f_k * g_k\| &\leq \|(f_j - f_k) * g_j\|_{\infty} + \|f_k * (g_j - g_k)\|_{\infty} \\ &\leq \|f_j - f_k\|_2 \|g_j\|_2 + \|f_k\|_2 \|g_j - g_k\|_2 \\ &\leq M \left(\|f_j - f_k\|_2 + \|g_j - g_k\|_2\right) \end{aligned}$$

where  $M \ge ||f_j||_2$  and  $M \ge ||g_k||_2$ , since the sequences converge in  $L^2(\mathbb{T})$ . Since  $C(\mathbb{T})$  is complete, the sequence  $(f_k * g_k)$  converges uniformly to a continuous function  $f * g \in C(\mathbb{T})$ , and f \* g satisfies the inequality.

Theorem 7.18. *Convolution Theorem* Theorem 7.5 on page 154 and Notes 1/10/11

If  $f, g \in L^2(\mathbb{T})$ , then

(Book) 
$$\widehat{(f * g)}_n = \sqrt{2\pi} \hat{f}_n \hat{g}_n$$
  
(Notes)  $\widehat{(f * g)}_n = \hat{f}_n \hat{g}_n$ 

#### Proof Outline

Compute  $(f * g)_n$ , using Fubini's Theorem to change the order of integration.

**Remark 7.19.** *Alternative bases for*  $L^2$  page 155 and Notes 1/7/11

The non-normalized orthogonal basis:

$$\{e^{inx}\}$$
$$\hat{f}_n = \frac{1}{2\pi} \int_{\mathbb{T}} f(x) e^{-inx} dx$$

Sines and Cosines:

$$\{1, \cos(nx), \sin(nx) \mid n = 1, 2, 3, \ldots\}$$

$$f(x) = \frac{1}{2}a_0 + \sum_{n=1}^{\infty} \left( a_n \cos(nx) + b_n \sin(nx) \right)$$
$$a_0 = \frac{1}{\pi} \int_{\mathbb{T}} f(x) \, dx \qquad a_n = \frac{1}{\pi} \int_{\mathbb{T}} f(x) \cos(nx) \, dx \qquad b_n = \frac{1}{\pi} \int_{\mathbb{T}} f(x) \sin(nx) \, dx$$

### 7.2 $L^1$ Functions

**Remark 7.20.**  $L^1$  *Functions* Notes 1/7/11

 $L^1(\mathbb{T})$  is the space of periodic functions  $f:\mathbb{T}\to\mathbb{C}$  such that

$$\|f\|_{L^1} = \int_{\mathbb{T}} |f(x)| \, dx < \infty$$

Note that  $L^1(\mathbb{T})$  is a Banach space but not a Hilbert space. We can define the Fourier coefficients of f as

$$c_n = \int_{\mathbb{T}} f(x) e^{-inx} \, dx$$

Note that  $|c_n| \leq \int |f(x)| dx$ . We can write the trigonometric polynomial approximation of f as

$$f(x) \sim \sum_{n \in \mathbb{Z}} c_n e^{inx}$$

However, this does not necessarily converge to f.

Lemma 7.21. *Riemann-Lebesgue Lemma* Notes 1/7/11 and 1/10/11

If  $f \in L^1(\mathbb{T})$  has Fourier coefficients  $c_n$ , then  $c_n \to 0$  as  $|n| \to \infty$ .

Proof Outline (1/7/11)

- Prove for smooth functions (use Integration By Parts)
- Approximate non-smooth functions with smooth functions

#### Proof Outline (1/10/11)

- Fix  $\epsilon > 0$
- The trigonometric polynomials are dense in  $L^1(\mathbb{T})$ , so we can pick a trigonometric polynomial p such that  $||f p||_{L^1} < \epsilon$
- If deg p = N and n > N, then

$$\begin{split} |\hat{f}(n)| &= \frac{1}{2\pi} \left| \int f e^{-inx} \, dx \right| \\ &= \frac{1}{2\pi} \left| \int (f-p) e^{-inx} \, dx \right| \qquad \text{Note: } \int p e^{-inx} \, dx = 0 \,\,\forall \, n > N \text{ by orthogonality} \\ &\leq \frac{1}{2\pi} \|f-p\|_{L^1} \\ &\leq \frac{\epsilon}{2\pi} < \epsilon \end{split}$$

**Definition 7.22.** Fourier Transform for  $L^1(\mathbb{T})$ Notes 1/10/11

The Fourier Transform  $\mathcal{F}: f \to \hat{f}, \ \mathcal{F}: L^1(\mathbb{T}) \to C_0(\mathbb{Z})$ 

$$C_0(\mathbb{Z}) = \left\{ (c_n)_{n \in \mathbb{Z}} \mid c_n \to 0 \text{ as } |n| \to \infty \right\}$$
$$\|(c_n)\|_{\infty} = \max_{n \in \mathbb{Z}} |c_n|$$

 $\mathcal{F}$  is a bounded linear map, with  $\|\mathcal{F}f\|_{\infty} \leq \|f\|_{L^1}$ Note:  $\mathcal{F}$  is not onto.

Example 7.23.  $\mathcal{F}$  is not onto Notes 1/10/11

There is no function whose Fourier coefficients are

$$\hat{f}(n) = \frac{i \operatorname{sgn}(n)}{\log |n|} \qquad |n| \ge 2$$

#### 7.3 Kernels and Summability Methods

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Definition 7.24. Dirichlet Kernel
Notes 1/10/11 and FA 44
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The *Dirichlet kernel* is

$$D_N(x) = \frac{1}{2\pi} \sum_{|n| \le N} e^{inx} = \frac{1}{2\pi} \left[ \frac{\sin\left((N + \frac{1}{2})x\right)}{\sin(\frac{x}{2})} \right] \quad x \ne 0$$
$$D_N(0) = \frac{1}{2\pi} (2N + 1)$$

(See the Kernel Overview.)

Derivation of the Dirichlet Kernel

Suppose  $f \in L^1(\mathbb{T}), \ f(x) \sim \sum \hat{f}_n e^{inx}$ . Define the N<sup>th</sup> partial sum of the Fourier series of f as

$$S_N(f)(x) = \sum_{|n| \le N} \hat{f}_n e^{inx}$$
  
=  $\frac{1}{2\pi} \sum_{|n| \le N} \left( \int f(y) e^{-iny} \, dy \right) e^{inx}$   
=  $\frac{1}{2\pi} \int \left( \sum_{|n| \le N} e^{in(x-y)} \right) f(y) \, dy$   
=  $\int D_N(x-y) f(y) \, dy = D_N * f$ 

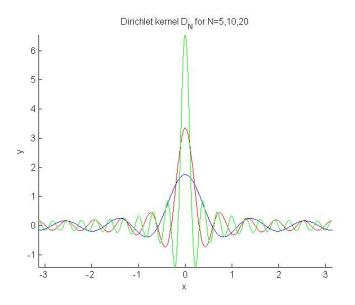


Figure 1: Dirichlet kernels.

#### **Example 7.25.** $D_N$ is not an approximate identity Notes 1/12/11

The Dirichlet kernel is not an approximate identity.

(a) 
$$\int D_N dx = \int \left(\frac{1}{2\pi} \sum e^{inx}\right) dx = \frac{1}{2\pi} \cdot 2\pi = 1$$
  
(b) 
$$\int \frac{4}{\pi^2} \sum_{k=1}^N \frac{1}{k} \le |D_N| dx \le \frac{4}{\pi^2} \left(\sum_{k=1}^N \frac{1}{k}\right) + 2 + \frac{\pi}{4}$$
  
As  $N \to \infty$ ,  $\int |D_N| dx = \frac{4}{\pi} \log N + O(1) \to \infty$  as  $N \to \infty$   
(c) For  $\delta > 0$ ,  $\lim_{N \to \infty} \int_{|x| > \delta} |D_N| dx \ne 0$ 

Thus, we can't conclude that if  $f \in C(\mathbb{T})$  or  $f \in L^1(\mathbb{T})$  then  $D_N * f \to f$  uniformly

#### Theorem 7.26. *Absolute Convergence* HW 3 Problem 2 and FA page 41

If  $f \in C(\mathbb{T})$  and its Fourier series is absolutely convergent,  $\sum_{n \in \mathbb{Z}} |\hat{f}(n)| < \infty$ , then the Fourier series converges uniformly to f.

Let  $\mathcal{A}(\mathbb{T})$  denote the space of integrable functions whose Fourier coefficients are absolutely convergent. That is,  $f \in \mathcal{F}(\mathbb{T})$  if  $\sum_{n \in \mathbb{Z}} |\hat{f}(n)| < \infty$ . If  $f \in \mathcal{A}(\mathbb{T})$ , then  $f \in C(\mathbb{T})$ .

**Definition 7.27.** Summability Method: Cesáro Summation Notes 1/12/11 and FA 52

The  $N^{\text{th}}$  Cesáro sum of a series is the average of the first N partial sums in the series:

$$\sigma_N = \frac{s_0 + s_1 + \ldots + s_{N-1}}{N}$$

**Example 7.28.** Cesáro Summation Example Notes 1/12/11

Consider the series  $\sum_{n=1}^{\infty} (-1)^n = 1 - 1 + 1 - 1 + 1 \dots$  Then the *n*th partial sum is

$$S_N = \begin{cases} 1 & N \text{ odd} \\ 0 & N \text{ even} \end{cases}$$

Consider the averages of partial sums:

$$\sigma_N = \frac{S_1 + \dots + S_N}{N}$$
$$\sigma_N = \begin{cases} \frac{1}{2} & N \text{ even} \\ \frac{\frac{1}{2}(N+1)}{N} = \frac{1}{2} + \frac{1}{2N} & N \text{ even} \end{cases} \rightarrow \frac{1}{2} \text{ as } N \rightarrow \infty$$

Thus, 
$$\sum_{n=1}^{\infty} (-1)^n = \frac{1}{2}$$
 (C).

## Theorem 7.29.

Notes 1/14/11

Cesáro summation is *regular*, meaning that if  $\sum a_n = s$  then  $\sum a_n = s$  (C).

Definition 7.30. Fejér Kernel Notes 1/12/11

The Fejér Kernel is:

$$K_N(x) = \frac{1}{2\pi} \sum_{|n| \le N} \left( 1 - \frac{|n|}{N+1} \right) e^{inx}$$
$$K_N(x) = \frac{1}{2\pi(N+1)} \left[ \frac{\sin\left(\frac{(N+1)x}{2}\right)}{\sin\left(\frac{x}{2}\right)} \right]^2$$

(See the Kernel Overview.)

 $\underline{Proof}$  (that the two forms are equivalent)

• Consider

$$\left[\frac{1}{2}\left(e^{ix} + e^{-ix}\right) - 1\right]K_N(x) = \frac{1}{2\pi N}\left(\frac{1}{2}e^{i(N+1)x} + \frac{1}{2}e^{-i(N+1)x} - 1\right)$$

• Use the fact that

$$\left(\sin\frac{x}{2}\right)^2 = -\frac{1}{4}\left(e^{ix} - 2 + e^{-ix}\right)$$

 $\frac{\text{Derivation of the Fejér Kernel}}{\text{Form the } N^{\text{th}} \text{ Cesáro mean of the Fourier series:}}$ 

$$\sigma_N(f)(x) = \frac{S_0 f + S_1 f + \dots + S_N f}{N+1}$$
$$= \frac{1}{2\pi} \sum_{|n| \le N} \left(1 - \frac{|n|}{N+1}\right) \hat{f}(n) e^{inx}$$
$$= K_N * f$$

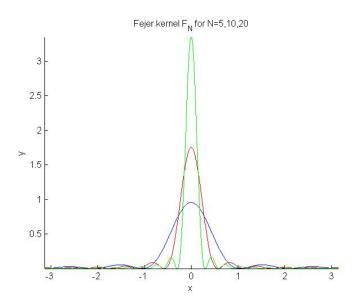


Figure 2: Fejér kernels.

#### **Theorem 7.31.** Notes 1/12/11

 $K_N$  is an approximate identity. If  $f \in C(\mathbb{T})$ , then  $\sigma_N f = K_N * f \to f$  uniformly and if  $f \in L^p(\mathbb{T})$ , then  $\sigma_N f = K_N * f \to f$  in  $L^p(\mathbb{T})$ .

# Corollary 7.32.

1/12/11

Suppose  $f, g \in L^1(\mathbb{T})$  and  $\hat{f} = \hat{g}$ . Then f = g.

Proof

- Set h = f g
- Then  $\hat{h}(n) = 0$
- $K_N * h \to h$  in  $L^1$
- $K_N * h = 0 \forall N$ , so  $h = 0 \Rightarrow f = g$

Note: we could have used the original approximate identity for this proof.

**Definition 7.33.** Summability Method: Abel Summation Notes 1/14/11

$$S = \sum_{n=0}^{\infty} a_n$$
$$S = \lim_{r \to 1^-} \sum_{n=0}^{\infty} a_n r^n \quad (A)$$

**Theorem 7.34.** Notes 1/14/11

Abel summation is regular.

 $\underline{\text{Proof}}$ 

• We will use summation by parts. Suppose  $S = \sum_{n=0}^{\infty} a_n$ ,  $S_n = \sum_{k=0}^{n} a_k$ ,  $S_n \to S$  as  $n \to \infty$ 

$$\begin{split} \sum_{n=0}^{\infty} a_n r^n &= a_0 + \sum_{n=1}^{\infty} (S_n - S_{n-1}) r^n & (\text{Since } a_n = S_n - S_{n-1}) \\ &= a_0 + \sum_{n=1}^{\infty} (S_n - S_n r^{n+1}) & (\text{re-index}) \\ &= a_0 + (1-r) \sum_{n=1}^{\infty} (S_n r^n) - S_0 r \\ &= (1-r) \sum_{n=0}^{\infty} S_n r^n & (S_0 = a_0) \\ \left| \sum_{n=0}^{\infty} (a_n r^n) - s \right| &= (1-r) \left| \sum_{n=0}^{\infty} (S_n - S) r^n \right| \le (1-r) \sum_{n=0}^{\infty} |S_n - S| r^n & 1 = (1-r) \sum_{n=0}^{\infty} r^n \\ &S = (1-r) \sum_{n=0}^{\infty} Sr^n \end{split}$$

• Fix  $\epsilon > 0$ . Choose N such that  $|S_n - S| < \epsilon/2$  for n > N. Then

$$\left|\sum_{n=0}^{\infty} a_n r^n - S\right| < (1-r) \sum_{n=0}^{N} |S_n - S| r^n + \frac{\epsilon}{2} \underbrace{(1-r) \sum_{n=N+1}^{\infty} r^n}_{\leq 1}$$

- Choose  $(1-r) < \delta$ , where  $\delta \sum_{n=0}^{N} |S_n S| < \epsilon/2$
- $n > N \Rightarrow \left|\sum_{n=0}^{\infty} a_n r^n S\right| < \epsilon/2 + \epsilon/2 = \epsilon$

#### Theorem 7.35. *Tauber & Littlwood* Notes 1/14/11

Suppose that  $\lim_{r \to 1^-} \sum_{n=0}^{\infty} a_n r^n$  exists and  $na_n = O(1)$  as  $n \to \infty$ . (i.e. there is an M such that  $|na_n| \leq M \forall n$ .) Then  $\sum a_n$  exists (and is equal to the limit).

# Definition 7.36. Poisson Kernel

Notes 1/14/11

Identify  $\mathbb T$  as the unit circle in  $\mathbb C,$  i.e.

$$\mathbb{T} = \left\{ z \in \mathbb{C} \mid |z| = 1 \right\} \Leftrightarrow z = e^{i\theta}$$
$$f(\theta) \sim \sum_{n \in \mathbb{Z}} \hat{f}(n) e^{in\theta}$$
$$\hat{f}(n) = \frac{1}{2\pi} \int_0^{2\pi} f(\theta) e^{-in\theta} d\theta$$
$$f_r(\theta) = \sum_{n \in \mathbb{Z}} \hat{f}(n) r^{|n|} e^{in\theta}$$
$$= P_r * f(\theta)$$

The Poisson kernel is

$$P_r(\theta) = \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{in\theta}, \quad 0 < r < 1$$
$$P_r(\theta) = \frac{1}{2\pi} \left[ \frac{1 - r^2}{1 - 2r \cos \theta + r^2} \right]$$
$$P_r(0) = \frac{1}{2\pi} \frac{1 - r^2}{(1 - r)^2}$$

(See the Kernel Overview.)

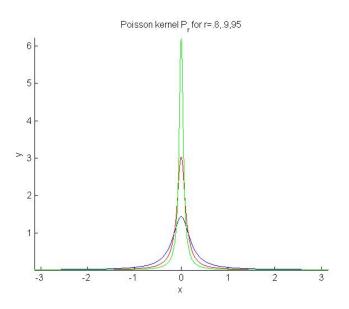


Figure 3: Poisson kernels.

**Remark 7.37.** *Properties of the Poisson Kernel* Notes 1/14/11

- The Poisson kernel is not a trigonometric polynomial
- The Poisson kernel satisfies:
  - (a)  $\int P_r(\theta) d\theta = 1$
  - (b)  $P_r \ge 0$
  - (c)  $P_r(\theta) \to 0$  uniformly as  $r \to 1^-$  on  $\delta < |\theta| < \pi$

# Theorem 7.38. Notes 1/1//11

Notes 1/14/11

 $P_r$  is an approximate identity as  $r \to 1^-$ .

**Corollary 7.39.** Notes 1/14/11

> If  $f \in L^p(\mathbb{T}), 1 \leq p < \infty$ , then  $P_r * f \to f$  as  $r \to 1^-$ . If  $f \in C(\mathbb{T})$ , then  $P_r * f \to f$  uniformly.

#### Remark 7.40. Kernel Overview

**Dirichlet** 

• Equations:

$$- D_N(x) = \frac{1}{2\pi} \sum_{|n| \le N} e^{inx} - D_N(x) = \frac{1}{2\pi} \left[ \frac{\sin\left((N + \frac{1}{2})x\right)}{\sin\left(\frac{x}{2}\right)} \right], \quad x \ne 0 - D_N(0) = \frac{1}{2\pi} (2N + 1)$$

- Summability Method: Standard
- Approximate Identity: No

Fejér

• Equations:

$$- K_N(x) = \frac{1}{2\pi} \sum_{|n| \le N} \left( 1 - \frac{|n|}{N+1} \right) e^{inx}$$
$$- K_N(x) = \frac{1}{2\pi(N+1)} \left[ \frac{\sin\left(\frac{(N+1)x}{2}\right)}{\sin\left(\frac{x}{2}\right)} \right]^2$$

- Summability Method: Cesáro
- Approximate Identity: Yes

#### <u>Poisson</u>

• Equations:

$$-P_r(\theta) = \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{in\theta}, \quad 0 < r < 1$$
$$-P_r(\theta) = \frac{1}{2\pi} \left[ \frac{1 - r^2}{1 - 2r \cos \theta + r^2} \right]$$
$$-P_r(0) = \frac{1}{2\pi} \frac{1 - r^2}{(1 - r)^2}$$

- Summability Method: Abel
- Approximate Identity: Yes, as  $r \to 1^-$

### 7.4 Harmonic Functions

#### **Definition 7.41.** *Harmonic* Notes 1/19/11

Let  $\Omega \subset \mathbb{R}^n$  be an open set.  $u: \Omega \to \mathbb{R}$  is *harmonic* on  $\Omega$  if  $\Delta u = 0$  in  $\Omega$ .

Recall:  $\Delta = \frac{\partial^2}{\partial x_1^2} + \frac{\partial^2}{\partial x_2^2} + \ldots + \frac{\partial^2}{\partial x_n^2}$ 

Remark 7.42. Harmonic & Analytic Functions Notes 1/19/11

There is a close connection in 2-D between harmonic and analytic (holomorphic) functions.

$$F: \Omega \to \mathbb{C}$$
$$F(z) = u(x, y) + iv(x, y)$$

where u, v satisfy the Cauchy-Riemann equations:

$$\left. \begin{array}{c} u_x = v_y \\ u_y = -v_x \end{array} \right\} \Rightarrow u_{xx} + v_{yy} = 0$$

**Example 7.43.**  $\Delta u = 0$  on the Complex Unit Disk Notes 1/19/11

Consider the Dirichlet problem on  $D = \{(x, y) \subset \mathbb{R}^2 \mid x^2 + y^2 < 1\}$ :

$$\Delta u = 0 \text{ in } D$$
$$u = f \text{ on } \partial D = \pi$$

Here  $f \in C(\partial D)$ . Want  $u \in C^2(D) \cap C(\overline{D})$ . Use separation of variables:

 $u(r,\theta) = F(r)G(\theta)$ 

We get that:

$$G(\theta) = e^{in\theta}$$
  

$$F(r) = Ar^{n} + Br^{-n} \quad n \neq 0$$
  

$$F(r) = A + B \ln r \quad n = 0$$

We want the solution to belong to  $C^2(D)$ , so we set

$$F(r) = r^{|n|}, \quad n \in \mathbb{Z}$$
$$\Rightarrow u(r, \theta) = \sum_{n \in \mathbb{Z}} c_n r^{|n|} e^{in\theta}$$

We want that:

$$u(1,\theta) = f(\theta) = \sum_{n \in \mathbb{Z}} c_n e^{in\theta}$$
$$\Rightarrow c_n = \hat{f}(n) = \frac{1}{2\pi} \int_0^{2\pi} f(\theta) e^{-in\theta} d\theta$$

Note that:

$$u(r,\theta) = \underbrace{(P_r * f)(\theta)}_{\text{Green's function}} = \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{in\theta}$$

#### **Remark 7.44.** Notes 1/19/11

 $P_r(\theta)$  is a  $C^{\infty}(D)$  function of  $r, \theta$  in  $0 \leq r < 1$ , and

$$\Delta P_r(\theta) = \frac{1}{r} \frac{\partial}{\partial r} \left( r \frac{\partial P_r}{\partial r} \right) + \frac{1}{r^2} \frac{\partial^2 P_r}{\partial \theta^2} = 0$$

# Theorem 7.45.

Notes 1/19/11

Suppose that  $f \in C(\partial D)$ . Then  $u(r, \theta) = (P_r * f)(\theta)$  is a solution of

$$\begin{cases} \Delta u = 0 & \text{in } D \\ u = f & \text{on } \partial D \end{cases}$$

Moreover,  $u \in C^{\infty}(D) \cap C(\overline{D})$ .

#### Proof

- $u(r,\theta) = \int_{\mathbb{T}} P_r(\theta \phi) f(\phi) \, d\phi$  (by Lebesgue Dominated Convergence Theorem)
- So  $u \in C^{\infty}(D)$ , and  $\Delta u = 0$
- Moreover,  $P_r * f \to f$  uniformly as  $r \to 1^-$
- So  $u \in C(\overline{D})$

## Theorem 7.46.

Notes 1/19/11

There is a unique solution  $u \in C^2(D) \cap C(\overline{D})$  of the Dirichlet problem. (Can be proved using the maximum principle and/or energy estimates.)

# Corollary 7.47.

Notes 1/19/11

Every harmonic function  $u \in C^2(D) \cap C(\overline{D})$  is smooth and has the mean value property:

$$u(r=0) = \frac{1}{2\pi} \int_0^{2\pi} f(\theta) \, d\theta$$

**Remark 7.48.** Background Info/Review Notes 1/21/11 • Function Spaces  $- \text{Let } 1 \leq p < \infty$ . If  $f \in L^p(\mathbb{T})$ , then  $f : \mathbb{T} \to \mathbb{C}$  and  $||f||_p = (\int_{\mathbb{T}} |f|^p dx)^{1/p} < \infty$ .  $- f = g \text{ in } L^p \text{ if } f = g \text{ a.e.}$   $- \text{ In } L^\infty, ||f||_\infty = \text{ess sup } |f(x)| = \inf_{\text{measure } N=0} \sup\{|f(x)| \mid x \in \mathbb{T} \setminus N\}$ • Sequence Spaces  $- \text{Let } 1 \leq q < \infty$ . If  $\hat{f} \in \ell^q(\mathbb{Z}), \ \hat{f} : \mathbb{Z} \to \mathbb{C}$ , then  $||\hat{f}||_q = \left(\sum_{n \in \mathbb{Z}} |\hat{f}(n)|^q\right)^{1/q} < \infty$   $- \text{ In } \ell^\infty, ||\hat{f}||_\infty = \sup_{n \in \mathbb{Z}} |\hat{f}(n)|$ • Question: When is  $\mathcal{F} : L^p(\mathbb{T}) \to \ell^q(\mathbb{Z}), \ f \mapsto \hat{f}$ , a bounded linear map?  $- \mathcal{F} : L^2 \to \ell^2$   $* ||\mathcal{F}f||_{\ell^2} = \frac{1}{\sqrt{2\pi}} ||f||_{L^2}$   $* \mathcal{F}$  is onto  $- \mathcal{F} : L^1 \to C_0 \subset \ell^\infty$ 

 $* \|\mathcal{F}f\|_{\ell^{\infty}} \leq \frac{1}{2\pi} \|f\|_{L^{1}}$ 

**Theorem 7.49.** *Hausdorff-Young Theorem/Inequality* Notes 1/21/11

Suppose  $1 \le p \le 2$  and  $2 \le p' \le \infty$  are Hölder conjugates  $(\frac{1}{p} + \frac{1}{p'} = 1)$ . Then  $\mathcal{F} : L^p(\mathbb{T}) \to \ell^{p'}(\mathbb{Z})$  is a bounded linear map, i.e.  $\|\hat{f}\|_{\ell^{p'}} \le C_p \|f\|_{L^p}$ .

## Remark 7.50.

Notes 1/21/11

- 1. Interpolation result (Riesz-Thorin Theorem)
- 2.  $\mathcal{F}$  is not onto if  $1 \leq p < 2$ .
  - Ex:  $p = 1, p' = \infty$ , then  $f \in L^1 \to \hat{f} \in C_0 \Rightarrow \text{not all of } \ell^{\infty}$
  - $\sum_{|n|\geq 2} \frac{i \operatorname{sgn}(n)}{\log n} e^{inx}$  is not the Fourier series of any  $L^1$  function
- 3. This result does not hold for 2
- 4. If  $f \in L^p$  (or even if  $f \in C$ ), one can't say much about the Fourier coefficients  $\hat{f}$  beyond the fact that  $f \in L^p$  so  $\hat{f} \in \ell^2$ 
  - Example:

$$\begin{split} f(x) &= \sum_{n=2}^{\infty} \frac{e^{in\log n}}{n^{1/2} (\log n)^2} e^{inx} \\ \hat{f}(n) &= \frac{e^{in\log n}}{n^{1/2} (\log n)^2} \\ \sum |\hat{f}(n)|^2 &= \sum \frac{1}{n (\log n)^4} < \infty \\ \hat{f} \in \ell^2 \text{ so } f \in L^2. \text{ Is } \hat{f} \in \ell^p \text{ for } p < 2, \text{ e.g. } p = 2 - \epsilon? \\ \sum |\hat{f}(n)|^{2-\epsilon} &= \sum \frac{1}{n^{1-\epsilon/2} (\log n)^{4-2\epsilon}} = \infty \end{split}$$

So  $\hat{f} \notin \ell^{p'}$  for any p' < 2

#### Fourier Series of Differentiable Functions (Section 7.2 in H&N) 7.6

**Definition 7.51.** Fourier Series Differentiation Notes 1/24/11

$$f(x) = \sum_{n \in \mathbb{Z}} c_n e^{inx}$$
$$f'(x) = \sum_{n \in \mathbb{Z}} inc_n e^{inx}$$
$$\mathcal{F} : \frac{d}{dx} \mapsto in$$

#### Proposition 7.52. Notes 1/24/11

If  $f \in C^1(\mathbb{T})$ , then  $\widehat{f'}(n) = in\widehat{f}(n)$ (Actually, it is sufficient that  $f \in L^1(\mathbb{T})$ .)

See Definition 7.56 and Proposition 11.21.

# Definition 7.53. Orders

Notes 1/24/11

If  $\phi, \psi : \mathbb{Z} \to \mathbb{C}$ , we say that

- $\phi = O(\psi)$  as  $|n| \to \infty$  if there exists C such that  $|\phi(n)| \le C|\psi(n)| \ \forall n \in \mathbb{Z}$
- $\phi = o(\psi)$  as  $|n| \to \infty$  if  $\lim_{|n| \to \infty} \left| \frac{\phi(n)}{\psi(n)} \right| = 0$

# Theorem 7.54.

Notes 1/24/11

If 
$$f \in C^1(\mathbb{T})$$
, then  $\hat{f}(n) = o(\frac{1}{n})$  as  $|n| \to \infty$ 

If  $f \in C^k(\mathbb{T})$ , where  $k \in \mathbb{N}$ , then  $\hat{f}(n) = o(\frac{1}{n^k})$  as  $|n| \to \infty$ 

Proof

- $\hat{f'}(n) = in\hat{f}(n)$  if  $f \in C^1$
- $\hat{f}(n) = \frac{1}{in}\hat{f}'(n), n \neq 0$ , and  $\hat{f}'(n) \to 0$  as  $|n| \to \infty$  by the Riemann-Lebesgue Lemma
- So  $\hat{f}(n) = o(\frac{1}{n})$  as  $|n| \to \infty$
- In general,  $\widehat{f}(n) = \frac{1}{(in)^k} \widehat{f^k}(n) = o(\frac{1}{n^k})$

#### Corollary 7.55. page 157 and Notes 1/24/11

If 
$$f \in C^{\infty}(\mathbb{T})$$
, then  $\lim_{|n| \to \infty} |n|^k \hat{f}(n) = 0 \ \forall \ k \in \mathbb{N}$ .

In other words, the Fourier coefficients of smooth functions form a rapidly decreasing sequence that decreases faster than any polynomial. Heuristically, a smooth function contains a small amount of high frequency components.

Compare to Theorem 11.18.

**Definition 7.56.** Weak  $L^2$ -derivatives (1) Notes 1/24/11

Suppose that  $f \in L^2(\mathbb{T})$  such that  $\sum_{n \in \mathbb{Z}} n^2 |\hat{f}(n)|^2 < \infty$ . Then we define the *weak*  $L^2$ -derivative  $g = f' \in L^2(\mathbb{T})$  by  $a(x) = \sum_{n \in \mathbb{Z}} in \hat{f}(n) e^{inx}$ 

$$g(x) = \sum_{n \in \mathbb{Z}} in\hat{f}(n)e^{inx}$$

See Proposition 7.52 and Proposition 11.21.

**Definition 7.57.** Sobolev Space (1) page 158 and Notes 1/24/11

$$H^{1}(\mathbb{T}) = \{f \in L^{2}(\mathbb{T}) \mid f' \in L^{2}(\mathbb{T})\}$$
$$\langle f, g \rangle_{H^{1}} = \int_{\mathbb{T}} (\overline{f}g + \overline{f'}g') \, dx = \sum_{n \in \mathbb{Z}} (1 + n^{2}) \overline{\widehat{f}(n)}g(n)$$
$$\|f\|_{H^{1}} = \left[\int_{\mathbb{T}} (|f|^{2} + |f'|^{2}) \, dx\right]^{1/2}$$

In other words,  $f \in H^1(\mathbb{T})$  iff f and its weak derivative f' (defined by integration by parts) belong to  $L^2(\mathbb{T})$ .

**Definition 7.58.** *Integration By Parts* Notes 1/24/11

For  $f, g \in H^1$ :

$$\int_{\mathbb{T}} \overline{f'g} \, dx = 2\pi \sum \overline{\widehat{f'(n)}} \hat{g}(n)$$
$$= 2\pi \sum \overline{in\hat{f}(n)} \hat{g}(n)$$
$$= -2\pi \sum \overline{\widehat{f(n)}} in\hat{g}(n)$$
$$= -2\pi \sum \overline{\widehat{f(n)}} g'(n)$$
$$= -\int_{\mathbb{T}} \overline{fg'} \, dx$$

# **Definition 7.59.** Weak Derivative (2) page 159 and Notes 1/24/11

A function  $g \in L^1(\mathbb{T})$  is the *weak derivative* of a function  $f \in L^1(\mathbb{T})$ , written g = f', if for every  $\phi \in C^{\infty}(\mathbb{T})$  we have

$$\int_{\mathbb{T}} f\phi' \, dx = -\int_{\mathbb{T}} g\phi \, dx$$

In other words, we are using integration by parts  $(\int_{\mathbb{T}} \overline{f'}g \, dx = -\int_{\mathbb{T}} \overline{f}g' \, dx)$ , to define f' pointwise a.e. We determine  $\hat{g}(n) \forall n$  by choosing  $\phi = e^{-inx}$ .

Compare to Distributional Derivative, Definition 11.10.

**Example 7.60.** Weak Derivative of f(x) = |x|Notes 1/26/11

 $f(x) = |x| \qquad -\pi < x < \pi$ 

 $f \in C(\mathbb{T})$ , but its standard derivative  $f' \notin C(\mathbb{T})$  because f'(0) and  $f'(\pi)$  don't exist. We shall see if g = f' (weak derivative) exists. We want:

$$\int g\phi \, dx = -\int f\phi' \, dx$$
$$= -\int_0^{\pi} x\phi' \, dx + \int_{-\pi}^0 x\phi' \, dx$$
$$= -x\phi \Big|_0^{\pi} + \int_0^{\pi} \phi \, dx + x\phi \Big|_{-\pi}^0 - \int_{-\pi}^0 \phi \, dx$$
$$= -\pi\phi(\pi) + \pi\phi(-\pi) + \int_{-\pi}^{\pi} \operatorname{sgn} x\phi \, dx$$

We conclude that  $\int f\phi x = -\int g\phi \, dx \,\,\forall \,\,\phi \in C^{\infty}(\mathbb{T})$  if  $g(x) = \operatorname{sgn} x$ .

**Example 7.61.** Weak Derivative of  $f(x) = \operatorname{sgn} x$ Notes 1/26/11

$$\int h\phi \, dx = -\int g\phi' \, dx$$
  
=  $-\int_{0}^{\pi} \phi' \, dx + \int_{-\pi}^{0} \phi' \, dx$   
=  $-\left[\phi(\pi) - \phi(0)\right] + \left[\phi(0) - \phi(-\pi)\right]$   
=  $2\left[\phi(0) - \phi(\pi)\right]$ 

There is no such  $h \in L^1$ . To see this, take  $\phi = \frac{1}{2\pi} e^{-inx} \in C^{\infty}(\mathbb{T})$ .

$$\hat{h}(n) = \frac{1}{\pi} [1 - e^{in\pi}] = \begin{cases} \frac{2}{\pi} & n \text{ odd} \\ 0 & n \text{ even} \end{cases}$$

This contradicts the Riemann-Lebesge Lemma, and therefore there is no such  $h \in L^1$ .

Proposition 7.62. Notes 1/26/11

f is weakly differentiable with  $f \in L^1$  iff it is absolutely continuous.

**Definition 7.63.** *Absolutely Continuous* http://en.wikipedia.org/wiki/Absolute\_continuity#Absolute\_continuity\_of\_functions

f is absolutely continuous if it has a derivative f' a.e., the derivative is Lebesgue integrable, and

$$f(x) = f(a) + \int_{a}^{x} f'(t) dt$$

**Theorem 7.64.** Notes 1/26/11

If f is weakly differentiable with weak derivative  $g = f' \in L^1(\mathbb{T})$ , then

$$\hat{g}(n) = inf(n)$$

Proof

$$\hat{g}(n) = \frac{1}{2\pi} \int g(x) e^{-inx} \, dx = -\frac{1}{2\pi} \int f(x) e^{-inx} \, dx = in\hat{f}(n)$$

Proposition 7.65. Notes 1/26/11

A function  $f\in L^2(\mathbb{T})$  has a weak derivative  $g\in L^2(\mathbb{T})$  iff

$$\sum_{n \in \mathbb{Z}} n^2 |\hat{f}(n)|^2 < \infty$$

and then

$$g(x) = \sum_{n \in \mathbb{Z}} in \hat{f}(n) e^{inx}$$

Definition 7.66. Sobolev Space (2) Notes 1/26/11

The Sobolev space  $W^{1,p}(\mathbb{T})$ ,  $1 \leq p \leq \infty$ , consists of all functions  $f : \mathbb{T} \to \mathbb{C}$  s.t.  $f \in L^p(\mathbb{T})$ ,  $f' \in L^p(\mathbb{T})$ . If p = 2, we write  $W^{1,2}(\mathbb{T}) = H^1(\mathbb{T})$  (where the H is because it is a Hilbert space).

A function  $f \in H^1(\mathbb{T})$  iff

$$\sum_{n\in\mathbb{Z}}(1+n^2)|\widehat{f}(n)|^2<\infty$$

and

$$\begin{split} \|f\|_{H^1} &= \left(\int |f|^2 \, dx + \int |f'|^2 \, dx\right)^{1/2} \\ &= \left(\|f\|_{L^2}^2 + \|f'\|_{L^2}^2\right)^{1/2} \\ &= \left(2\pi \sum_{n \in \mathbb{Z}} (1+n^2) |\hat{f}(n)|^2\right)^{1/2} \end{split}$$

**Theorem 7.67.** Sobolev Embedding Theorem Notes 1/26/11

If 
$$f \in H^1(\mathbb{T})$$
 then  $f \in C(\mathbb{T})$  and

$$\|f\|_{\infty} \le C \|f\|_{H^1}$$

 $J: H^1 \to C \text{ (Embedding)}, \ f \mapsto f.$ 

Proof

$$\begin{split} \sum_{n \in \mathbb{Z}} |\hat{f}(n)| &= \sum_{n \in \mathbb{Z}} \frac{1}{(1+n^2)^{1/2}} (1+n^2)^{1/2} |\hat{f}(n)| \\ &\leq \left( \sum_{n \in \mathbb{Z}} \frac{1}{(1+n^2)^{1/2}} \right) \left( \sum_{n \in \mathbb{Z}} (1+n^2) |\hat{f}(n)| \right) \\ &\leq C \|f\|_{H^1} \end{split}$$

It follows that  $f \in C(\mathbb{T})$  because the Fourier series converges uniformly to f (see Theorem 7.26) and

$$\|f\|_{\infty} \le \sum_{n \in \mathbb{Z}} |\hat{f}(n)| \le C \|f\|_{H^{\frac{1}{2}}}$$

#### 7.7 Chapter Summary

This chapter explores the spaces  $L^p(\mathbb{T})$ ,  $p \in [1, \infty)$ , with special attention given to the Hilbert space  $L^2(\mathbb{T})$ . These spaces are the completion of  $C(\mathbb{T})$  with respect to the  $L^p$ -norm; thus,  $C(\mathbb{T})$  is dense in  $L^p(\mathbb{T})$  for  $p \in [1, \infty)$ . Since  $\mathbb{T}$  has finite Lebesgue measure, we can use Hölder's Inequality to show that for p > q,  $\|\cdot\|_p \geq \|\cdot\|_q$ , which implies that  $L^p(\mathbb{T}) \subset L^q(\mathbb{T})$ . We define the *convolution* of two functions and what it means for a family of functions to be an *approximate identity*, and we use these tools to prove the *Weierstrass Approximation Theorem*, which says that the trigonometric polynomials are dense in  $C(\mathbb{T})$  with respect to the uniform norm. Since uniform convergence implies  $L^2$  convergence, it follows that the functions  $e_n(x) = \frac{1}{\sqrt{2\pi}}e^{inx}$  form an orthonormal basis for  $L^2(\mathbb{T})$ . Thus, for all  $f \in L^2(\mathbb{T})$ , we have that

$$f(x) = \sum_{n = -\infty}^{\infty} \hat{f}_n e^{inx},$$

where the equality is in the  $L^2$  sense. A result from Carleson tells us that the Fourier series of f converges pointwise to f a.e.

Next we explore some properties of Fourier series and Fourier coefficients. Let  $f, g \in L^2(\mathbb{T})$ . We use the density of  $C(\mathbb{T})$  in  $L^2(\mathbb{T})$  to prove the *Convolution Theorem*, which allows us to express the Fourier coefficients of f \* g in terms of those of f and g:  $(f * g)_n = \sqrt{2\pi} \hat{f}_n \hat{g}_n$ . *Parseval's Theorem* allows us to compute  $\langle f, g \rangle$  using the Fourier coefficients of f and g:  $\langle f, g \rangle = \sum_{n=-\infty}^{\infty} \hat{f}_n \hat{g}_n$ .

Now we examine the Fourier series of differentiable functions. Using integration by parts, we show that

$$\hat{f}'_n = in\hat{f}_n$$

This gives us the concept of a weak derivative, since the derivative of f may not be continuous; e.g. f(x) = |x|. We define the Sobolev space  $H^k(\mathbb{T})$  as the space of  $L^2(\mathbb{T})$  functions with k weak derivatives. And since the boundary terms on  $\mathbb{T}$  vanish, we have that  $\langle f', g \rangle = -\langle f, g' \rangle$  for  $f, g \in H^1(\mathbb{T})$ . Thus, we may define the weak derivative of a function using integration by parts:  $g \in L^1(\mathbb{T})$  is the weak derivative of  $f \in L^1(\mathbb{T})$  if

$$\int_{\mathbb{T}} f \phi' \, dx = - \int_{\mathbb{T}} g \phi \, dx \qquad \forall \ \phi \in C^{\infty}(\mathbb{T}).$$

Finally, we prove a special case of the Sobolev Embedding Theorem: if  $f \in H^k(\mathbb{T})$  for k > 1/2, then  $f \in C(\mathbb{T})$ .

In addition, Hunter briefly discussed  $L^1(\mathbb{T})$ . We can define the Fourier series of an  $L^1$  function, but we cannot guarantee that it converges to the function. Our main result is the *Riemann-Lebesgue Lemma*, which says that the Fourier coefficients of an  $L^1$  function decay to zero as  $n \to \infty$ . Hunter then discussed 3 kernels: the *Dirichlet kernel* (standard summation), *Fejér kernel* (*Cesáro summation*), and *Poisson kernel* (*Abel summation*). These kernels are related to the concept of approximate identities, and we convolve the kernels with a function f. He covered harmonic functions, and our main result is that we can use the Poisson kernel to solve the two-dimensional Laplace equation.

## 11 Distributions and the Fourier Transform

#### 11.1 Periodic Distributions

#### Definition 11.1. Test Functions

Notes 1/28/11 and http://en.wikipedia.org/wiki/Distribution\_%28mathematics%29 and Hunter's Notes page 51

We define our space of *test functions* as:

 $\mathcal{D}(\mathbb{T}) = C^{\infty}(\mathbb{T})$  with the following topology:

 $\varphi_n \to \varphi \in \mathcal{D}$  if  $\varphi_n^{(k)} \to \varphi^{(k)}$  uniformly for all  $k = 0, 1, 2, \ldots$  Note that this topology is not obtained from any norm, but rather it is derived.

Definition 11.2. Distribution

Notes 1/28/11 and Hunter's Notes page 51

A distribution is a continuous linear functional, T, that maps a set of test functions,  $\mathcal{D}(\mathbb{T})$ , onto the set of complex numbers. The space of distributions is denoted by  $\mathcal{D}'(\mathbb{T})$ . For  $T \in \mathcal{D}'(\mathbb{T})$ ,  $\varphi \in \mathcal{D}(\mathbb{T})$ , we write:

$$\langle T, \varphi \rangle = T(\varphi)$$

 $\mathcal{D}'(\mathbb{T})$  is the topological dual space of the distributions on  $\mathbb{T}$  (i.e.  $\mathcal{D}(\mathbb{T})$ ), with the topology defined as follows:  $T_n \to T$  in  $\mathcal{D}'$  if  $\langle T_n, \varphi \rangle \to \langle T, \varphi \rangle$  in  $\mathbb{C} \forall \varphi \in \mathcal{D}$ .

 $T: \mathcal{D}(\mathbb{T}) \to \mathbb{C}$ Linear:  $\langle T, \lambda \varphi + \mu \psi \rangle = \lambda \langle T, \varphi \rangle + \mu \langle T, \psi \rangle$ Continuous: If  $\varphi_n \to \varphi \in \mathcal{D}$ , then  $\langle T, \varphi_n \rangle \to \langle T, \varphi \rangle \in \mathbb{C}$ 

Compare Distributional Convergence,  $T_n \rightharpoonup T$  in  $\mathcal{D}'$  if  $\langle T_n, \varphi \rangle \rightarrow \langle T, \varphi \rangle$ , to Weak Convergence (Definition 8.41):  $x_n \rightharpoonup x$  if  $\langle x_n, y \rangle \rightarrow \langle x, y \rangle \quad \forall y \in \mathcal{H}$ .

**Definition 11.3.** *Seminorm* Notes 1/28/11

Our topology on  $\mathcal{D}$  is obtained from a countable family of *seminorms*:

$$\|\varphi\|_k = \sup_{x \in \mathbb{T}} |\varphi^{(k)}(x)|, \qquad k = 0, 1, 2, \dots$$

A seminorm has the same properties as a norm except that it may assign length zero to nonzero vectors.

Example 11.4. Seminorms Notes 1/28/11

$$d(\varphi,\psi) = \sum_{k=0}^{\infty} \frac{1}{2^k} \frac{\|\varphi - \psi\|_k}{1 + \|\varphi - \psi\|_k}$$

- This is not a norm because you can't pull out a constant
- This turns  $\mathcal{D}$  into a *Fréchet space* (a complete, metrizable topological vector space topology defined by a countable family of seminorms)
- We could instead use norms to define the topology on  $\mathcal{D}(\mathbb{T})$ :

$$\|\varphi\|_{C^k} = \sum_{j=0}^k \|\varphi\|_j$$

# Remark 11.5.

Notes 1/28/11

Note that the differentiation operator

$$D: \mathcal{D}(\mathbb{T}) \to \mathcal{D}(\mathbb{T}), \quad D(\varphi) = \varphi'$$

is continuous: if  $\varphi_n \to \varphi \in \mathcal{D}$ , then  $D\varphi_n \to D\varphi \in \mathcal{D}$ . This is because there are infinitely many semi-norms.

# **Example 11.6.** *Regular Distribution* page 292 and Notes 1/28/11

If  $f: \mathbb{T} \to \mathbb{C}$  is integrable,  $f \in L^1(\mathbb{T})$ , define

$$T_f: \mathcal{D}(\mathbb{T}) \to \mathbb{C}$$
  
 $T_f(\varphi) = \int_{\mathbb{T}} f\varphi \, dx$ 

 $|T_f(\varphi)| \leq \sup |\varphi| \cdot \int |f| \, dx < \infty$ , so  $T_f$  is well-defined. It is a distribution because it satisfies:

- 1. Linearity: (1)  $T_f(\varphi + \psi) = \int f(\varphi + \psi) dx = T_f(\varphi) + T_f(\psi)$ . (2)  $T_f(c\varphi) = cT_f(\varphi)$
- 2. Continuity: If  $\varphi_n \to 0$  in  $\mathcal{D}$ , then  $|T_f(\varphi_n)| \leq \sup |\varphi_n| ||f||_{L^1} \to 0$  as  $n \to \infty$ . So  $T_f(\varphi_n) \to 0$  and  $T_f$  is continuous.

We identify f with  $T_f$ . Thus,  $L^1(\mathbb{T}) \subset D^1(\mathbb{T})$ .

We call  $T_f$  a regular distribution. A regular distribution is a distribution that is given by the integration of a test function  $\varphi$  against a function f.

#### **Definition 11.7.** *Principal Value Distribution* page 293

A principal value distribution is a singular distribution, denoted by p.v. (1/x), and its action on a test function  $\varphi$  is given by

p.v. 
$$\frac{1}{x}(\varphi) = \lim_{\epsilon \to 0^+} \int_{|x| > \epsilon} \frac{\varphi(x)}{x} dx$$

#### Example 11.8. Notes 1/28/11

Consider the periodic  $\delta$ -function (actually a distribution, not a function).

$$\begin{aligned} \langle \delta, \varphi \rangle &= \varphi(0) \\ \langle \delta, \varphi + \psi \rangle &= (\varphi + \psi)(0) = \varphi(0) + \psi(0) = \langle \delta, \varphi \rangle + \langle \delta, \psi \rangle \\ \langle \delta, c\varphi \rangle &= c \langle \delta, \varphi \rangle \end{aligned}$$

 $\varphi_n \to 0$  implies  $\varphi_n(0) \to 0$ , and therefore  $\delta$  is a continuous linear functional.

 $\delta$  is not regular. Proof:

- Suppose  $\langle \delta, \varphi \rangle = \int f \varphi \, dx$  for some  $f \in L^1$ .
- Consider  $\varphi_n(x) = \left[\frac{1+\cos x}{2}\right]^n$
- $\langle \delta, \varphi_n \rangle = 1 \forall n$ , but  $\int f \varphi_n \, dx \to 0$  as  $n \to \infty$  by the Lebesge-Dominated Convergence Theorem if  $f \in L^1$
- Thus, there is no function  $f \in L^1$  such that  $\int f\varphi \, dx = \varphi(0)$

#### Example 11.9.

Notes 1/28/11

Let 
$$T_n = \begin{cases} \frac{1}{2}n & |x| \le \frac{1}{n} \\ 0 & \frac{1}{n} \le |x| \le \pi \end{cases}$$

Then  $\int_{-\pi}^{\pi} T_n \, dx = 1 \, \forall n$ . Claim:  $\langle T_n, \varphi \rangle = \frac{n}{2} \int_{1/n}^{1/n} \varphi(x) \to \varphi(0)$  as  $n \to \infty$ . Proof:

$$\left| \frac{n}{2} \int_{-1/n}^{1/n} \varphi(x) \, dx - \varphi(0) \right| = \frac{n}{2} \left| \int_{-1/n}^{1/n} \left[ \varphi(x) - \varphi(0) \right] \, dx \right|$$
$$\leq \frac{n}{2} \left[ \sup_{|x| \le 1/n} \left| \varphi(x) - \varphi(0) \right| \right] \cdot \frac{2}{n}$$
$$\leq \sup_{|x| \le 1/n} \left| \varphi(x) - \varphi(0) \right| \to 0 \text{ as } n \to \infty$$

**Definition 11.10.** *Distributional Derivative* page 295

Every distribution  $T \in \mathcal{D}'(\mathbb{T})$  has a distributional derivative  $T' \in \mathcal{D}(\mathbb{T})$  that is given by

$$\langle T', \phi \rangle = - \langle T, \phi' \rangle \qquad \forall \phi \in \mathcal{D}(\mathbb{T})$$

Compare to Weak Derivative (2), Definition 7.59.

**Definition 11.11.** Motivation for Distributional Derivatives Notes 1/31/11

Suppose  $f \in C^{\infty}$  is a smooth function. Consider  $T_{f'}$ :

$$\langle T_{f'}, \varphi \rangle = \int f' \varphi \, dx = - \int f \varphi' \, dx = - \langle T_f, \varphi' \rangle$$

Want:  $(T_{f'}) = (T_f)'$ 

This defines the *distributional derivative*.

- 1. Linearity:  $\langle T', a\varphi + b\psi \rangle = -\langle T, (a\varphi + b\psi)' \rangle = -\langle T, a\varphi' + b\psi' \rangle = -a \langle T, \varphi' \rangle b \langle T, \psi' \rangle = a \langle T', \varphi \rangle + b \langle T', \psi \rangle$
- 2. Continuity: Suppose  $\varphi_n \to \varphi$  in  $\mathcal{D}$ . Consider  $\langle T', \varphi \rangle$ .  $\langle T', \varphi_n \rangle = -\langle T, \varphi'_n \rangle \to -\langle T, \varphi' \rangle = \langle T', \varphi \rangle$ , because T is continuous on  $\mathcal{D}$  and  $D : \varphi \to \varphi'$  is continuous on  $\mathcal{D}$

**Example 11.12.** Notes 1/31/11

 $f(x) = |x|, \quad |x| \le \pi$  $f'(x) = \operatorname{sgn} x = g(x)$ 

Compute the distributional derivative of g:

$$\begin{split} \left\langle g',\varphi\right\rangle &= -\left\langle g,\varphi'\right\rangle \\ &= -\int_0^\pi \varphi'\,dx + \int_{-\pi}^0 \varphi'\,dx \\ &= -\left[\varphi(\pi) - \varphi(0)\right] + \left[\varphi(0) - \varphi(\pi)\right] \\ &= 2\varphi(0) - 2\varphi(\pi) \\ &= 2\left\langle \delta_0,\varphi\right\rangle - 2\left\langle \delta_\pi,\varphi\right\rangle \\ &= \left\langle 2\delta_0 - 2\delta_\pi,\varphi\right\rangle \end{split}$$

$$g' = 2\delta_0 - 2\delta_\pi$$
$$= 2(\delta - \tau_\pi \delta)$$

Where  $\tau_{\pi}$  means translation by  $\pi$  and  $\delta_a$  is the  $\delta$ -"function" supported at a:

 $\langle \delta_a, \varphi \rangle = \varphi(a)$ 

#### **Example 11.13.** Notes 1/31/11

Compute  $\delta'$ :

$$\langle \delta', \varphi \rangle = - \langle \delta, \varphi' \rangle = -\varphi'(0)$$

**Definition 11.14.** *Fourier Coefficients* Notes 1/31/11

If  $T \in \mathcal{D}'(\mathbb{T})$ , define  $\hat{T}(n) = \frac{1}{2\pi} \langle T, e^{-inx} \rangle$ .

#### **Example 11.15.** Notes 1/31/11

Compute the Fourier coefficients of  $\delta$ :

$$\begin{split} \hat{\delta}(n) &= \frac{1}{2\pi} \left< \delta, e^{-inx} \right> = \frac{1}{2\pi} e^0 = \frac{1}{2\pi} \\ \delta(x) &= \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} e^{inx} \end{split}$$

#### Remark 11.16. 1/31/11

There are 3 contexts in which to look at Fourier series:

- Continuous functions  $\Rightarrow$  converge uniformly
- $L^2$  functions  $\Rightarrow$  converge in  $L^2$
- Distribution functions  $\Rightarrow$  converge in the distributional sense

**Example 11.17.** Notes 1/31/11

$$P_r(x) = \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{inx}$$
  
Formally, as  $r \to 1^-$ ,  $P_r(x) \rightharpoonup \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} r^{|n|} e^{inx} = \delta(x)$ 

#### **Theorem 11.18.** Notes 1/31/11

 $\varphi \in \mathcal{D}$  iff  $(\hat{\varphi}(n))$  is rapidly decreasing, i.e.

$$|n|^k \hat{\varphi}(n) \to 0 \text{ as } n \to \infty \ \forall \ k \ge 0$$

and the Fourier series of  $\varphi$  converges to  $\varphi$  in  $\mathcal{D}$ .

Compare to Corollary 7.55.

#### Proof

- $\varphi \in C^k \Rightarrow |n|^k \hat{\varphi}(n) \to 0$  by the Riemann-Lebesgue Lemma, so if  $\varphi \in C^\infty$ , then the  $\hat{\varphi}(n)$  are rapidly decreasing
- Sobolev Embedding Theorem: If  $\hat{\varphi}(n)$  is rapidly decreasing, then  $\varphi \in H^k(\mathbb{T}) \; \forall \; k$  implies that

 $\sum (1+n^2) |\hat{\varphi}(n)|^2 < \infty$ 

- Hence,  $\varphi \in C^{k-1}(\mathbb{T}) \ \forall \ k$ . So  $\varphi \in C^{\infty}$ .
- Similarly,  $\sum_{|n| \le N} \hat{\varphi}(n) e^{inx} \to \varphi$  in  $H^k \forall k$ - So  $\sum_{|n| \le N} \hat{\varphi}(n) e^{inx} \to \varphi$  in  $C^{k-1} \forall k$ - So  $\sum_{|n| \le N} \hat{\varphi}(n) e^{inx}$  converges in  $\mathcal{D}$

**Definition 11.19.**  $S(\mathbb{Z})$ Notes 2/2/11

 $S(\mathbb{Z})$  is the space of rapidly decreasing sequences,  $(c_n)$ , such that

$$\lim_{n \to \infty} |n|^k c_n = 0 \quad \forall \ k = 0, 1, 2, \dots$$

#### **Remark 11.20.** Notes 2/2/11

$$\begin{aligned} \mathcal{F} &: C^{\infty}(\mathbb{T}) \to S(\mathbb{Z}) \\ \mathcal{F} &: \varphi \to (\hat{\varphi}(n)) \end{aligned}$$
  
If  $\varphi \in C^{\infty}(\mathbb{T})$ , then  $S_N \varphi = \sum_{|n| \le N} \hat{\varphi}(n) e^{inx} \to \varphi$  in  $\mathcal{D}$ .  
If  $T \in \mathcal{D}'(\mathbb{T})$ , then  $\hat{T}(n) = \frac{1}{2\pi} \langle T, e^{-inx} \rangle$ 

Proposition 11.21. Notes 2/2/11

$$\widehat{T'}(n) = in\widehat{T}(n)$$

See Proposition 7.52 and Definition 7.56.

Proof.

$$\hat{T}'(n) = \frac{1}{2\pi} \left\langle T', e^{-inx} \right\rangle = -\frac{1}{2\pi} \left\langle T, \left( e^{-inx} \right)' \right\rangle = in \cdot \frac{1}{2\pi} \left\langle T, e^{-inx} \right\rangle$$
$$= in\hat{T}(n)$$

#### **Definition 11.22.** *Slow Growth* Notes 2/2/11

A sequence  $(c_n)$  has slow growth if there exist k, M such that  $|c_n| \leq M(1+n^{2k})^{1/2} \forall n$ .

Equivalently,  $|c_n| \leq M |n|^k \ \forall \ n \neq 0.$ 

Lemma 11.23. Notes 2/2/11

If  $T \in \mathcal{D}'$ , then  $(\hat{T}(n))$  has slow growth.

*Proof.* If  $T \in \mathcal{D}'$  then T has some finite order k such that

$$|\langle T, \varphi \rangle| \le C \|\varphi\|_{C^k}$$

Then

$$|\hat{T}(n)| = |\langle T, e^{-inx} \rangle| \le C ||e^{-inx}||_{C^k} \le C(1+n^{2k})^{1/2}$$

**Example 11.24.** Weierstrass Nowhwere Differentiable Function Notes 2/2/11

$$f(x) = \sum_{n=1}^{\infty} \frac{1}{2^n} \cos(3^n x)$$

 $\sum \frac{1}{2^n} < \infty$ , so  $f \in \mathcal{A}(\mathbb{T})$ .

$$f'(x) = \sum_{n=1}^{\infty} \frac{3^n}{2^n} \sin(3^n x)$$

f is nowhere differentiable, although it does have a distributional derivative.

# Theorem 11.25.

Notes 2/2/11

If 
$$T \in \mathcal{D}'(\mathbb{T})$$
 and  $S_N T = \sum_{|n| \le N} \hat{T}(n) e^{inx} \in C^{\infty}(\mathbb{T})$ , then  $S_N T \rightharpoonup T$  in  $\mathcal{D}'$  as  $N \rightarrow \infty$ .  
Ex:  $\delta(x) = \frac{1}{2\pi} \sum_{n \in \mathbb{Z}} e^{inx}$ 

Proof.

$$\langle S_N T, \varphi \rangle = \left\langle \sum_{|n| \le N} \hat{T}(n) e^{inx}, \varphi \right\rangle = \sum_{|n| \le N} \left\langle \hat{T}(n) e^{-inx}, \varphi \right\rangle = \sum_{|n| \le N} \hat{T}(n) \int e^{inx} \varphi(x) \, dx$$

$$= 2\pi \sum_{|n| \le N} \hat{T}(n) \hat{\varphi}(-n) = 2\pi \sum_{|n| \le N} \left\langle T, e^{-inx} \right\rangle \cdot \frac{1}{2\pi} \hat{\varphi}(-n) = \left\langle T, \sum_{|n| \le N} \hat{\varphi}(-n) e^{-inx} \right\rangle$$

$$= \left\langle T, S_N \varphi \right\rangle \to \left\langle T, \varphi \right\rangle \text{ as } n \to \infty$$

So  $S_N T \to T$  as  $N \to \infty$ .

# Theorem 11.26.

Notes 2/2/11

If  $(c_n)$  is a sequence of slow growth,  $(c_n) \in S'(\mathbb{Z})$ , then there exists a distribution T such that  $\hat{T}(n) = c_n$ .

*Proof.* Define T by

$$\langle T, \varphi \rangle = 2\pi \sum_{n \in \mathbb{Z}} c_n \hat{\varphi}(-n)$$

Remark 11.27.

 Notes 2/2/11

 
$$\mathcal{F}: f \mapsto \hat{f}(n)$$
 $\mathcal{D}(\mathbb{T}) = C^{\infty}(\mathbb{T}) \leftrightarrow S(\mathbb{Z})$ 
 $C(\mathbb{T}) \supset \mathcal{A}(\mathbb{T}) \leftrightarrow \ell'(\mathbb{Z})$ 
 $L^2(\mathbb{T}) \leftrightarrow \ell^2(\mathbb{Z})$ 
 $L^1(\mathbb{T}) \to C_0(\mathbb{Z})$ 
 $\mathcal{D}'(\mathbb{T}) \leftrightarrow S'(\mathbb{Z})$ 

 •  $C^{\infty} \subset L^2(\mathbb{T}) \subset \mathcal{D}'(\mathbb{T})$ 

 •  $S(\mathbb{Z}) \subset \ell^2(\mathbb{Z}) \subset S'(\mathbb{Z})$ 

## 8 Bounded Linear Operators on a Hilbert Space

#### 8.1 Orthogonal Projections

**Definition 8.1.** *Direct Sum* page 187 and Notes 2/4/11

If M and N are subspaces of a linear space X such that every  $x \in X$  can be written uniquely as x = y + z with  $y \in M$  and  $z \in N$ , then we say that  $X = M \oplus N$  is the *direct sum* of M and N, and we call N a *complementary subspace* of M in X. The decomposition x = y + z is unique if and only if  $M \cap N = \{0\}$ .

**Definition 8.2.** Projection, Idempotent, Self-Adjoint page 187 & 188 and Notes 2/4/11

Given a direct sum decomposition,  $X = M \oplus N$ , define the projection  $P : X \to X$  onto M along N by

 $P(m+n) = m, \qquad m \in M, \quad n \in N$ 

All projections are linear and *idempotent*, meaning that  $P^2 = P$ , because

$$P^2(m+n) = P(m) = m$$

Theorem 8.3. 100 kN  $\pm$  0/4

page 188 and Notes 2/4/11

Any linear map  $P: X \to X$  with  $P^2 = P$  is a projection. Specifically, it is the projection onto ran P along ker P.

Proof.

- x = P(x) + (x P(x))
- $P^2(x) = P(x) \implies P(x) \in \operatorname{ran} P$
- $P(x P(x)) = Px P^2x = Px Px = 0 \implies x P(x) \in \ker P$
- Suppose  $x \in \ker P \cap \operatorname{ran} P$ 
  - $-x \in \operatorname{ran} P \quad \Rightarrow \quad x = Py$
  - $-x \in \ker P \quad \Rightarrow \quad 0 = Px = P^2y = Py = x = 0$
  - Thus, x = 0, and ker  $P \cap \operatorname{ran} P = \{0\}$
- Thus,  $X = \operatorname{ran} P \oplus \ker P$

#### **Remark 8.4.** *Bounded Projections* Notes 2/4/11

**Question:** Given a projection  $P: X \to X, X$  a Banach space, when can we say that P is bounded?

**Answer:** We need ran P closed and complemented by a closed subspace  $N = \ker P$ 

Note: The kernel of a bounded operator is always closed; the range need not be.

#### **Definition 8.5.** Orthogonal Projections, Self-Adjoint Notes 2/4/11 and 2/7/11

Let  $\mathcal{H}$  be a Hilbert space and let  $M \subset \mathcal{H}$  be a closed linear subspace. Then by the Projection Theorem,

 $\mathcal{H} = M \oplus M^{\perp}, \qquad M^{\perp} = \{ y \in \mathcal{H} \mid y \perp m \; \forall \; m \in M \}$ 

We define the orthogonal projection  $P: \mathcal{H} \to \mathcal{H}$  onto M along  $M^{\perp}$ .

An orthogonal projection P on a Hilbert space  $\mathcal{H}$  is

- Idempotent:  $P^2 = P$
- Self-Adjoint:  $\langle x, Py \rangle = \langle Px, y \rangle$

*Proof.* To see that a projection P on a Hilbert space  $\mathcal{H}$  is self-adjoint, let

 $x = m + n, \qquad y = p + q, \qquad \text{where} \qquad m, p \in M, \qquad n, q \in N$ 

Compute:

$$\begin{aligned} \langle x, Py \rangle &= \langle m+n, p \rangle = \langle m, p \rangle + \langle n, p \rangle = \langle m, p \rangle \\ \langle Px, y \rangle &= \langle m, p+q \rangle = \langle m, p \rangle + \langle m, q \rangle = \langle m, p \rangle \end{aligned}$$

Lemma 8.6.

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If P is a nonzero othogonal projection then ||P|| = 1

Proof.

$$||Px||^{2} = \langle Px, Px \rangle = \langle x, P^{2}x \rangle = \langle x, Px \rangle \le ||x|| ||Px||$$

Either ||Px|| = 0 or  $||Px|| \le ||x||$ . Since  $||Px|| \ne 0 \forall x$ , it must be the case that  $||Px|| \le ||x||$ . Then

$$||P|| = \sup \frac{||Px||}{||x||} \le 1$$

If  $P \neq 0$ , then there exists  $y \in \mathcal{H}$  such that  $Py \neq 0$ . Setting x = Py in the previous equation yields

$$||P|| \ge \frac{||P \cdot Px||}{||Px||} = 1$$

# Theorem 8.7.

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If P is an orthogonal projection, then  $\mathcal{H} = M \oplus M^{\perp} = \operatorname{ran} P \oplus \ker P$ , where  $M = \operatorname{ran} P$  and  $M^{\perp} = \ker P$  are closed subspaces. Conversely, if M is any closed subspace of  $\mathcal{H}$ , then there exists an orthogonal projection with  $M = \operatorname{ran} P$  and  $M^{\perp} = \ker P$ .

**Example 8.8.** Even & Odd Functions page 189 and Notes 2/7/11

Let  $\mathcal{H} = L^2(\mathbb{R})$  and let

M = space of even functions, f(-x) = f(x)N = space of odd functions, f(-x) = -f(x)

 $M \perp N$ , since  $\int \overline{f}g \, dx = 0$  for f odd, g even. Define

- Even Projection:  $P: \mathcal{H} \to \mathcal{H}$  onto  $M, Pf(x) = \frac{1}{2}[f(x) + f(-x)]$
- Odd Projection:  $Q : \mathcal{H} \to \mathcal{H}$  onto  $N, Qf(x) = \frac{1}{2}[f(x) f(-x)]$ - Note: Q = I - P

Check that P is self-adjoint:

$$\langle Pf,g\rangle = \int_{\mathbb{R}} \frac{1}{2} \overline{[f(x) + f(-x)]} g(x) \, dx = \int_{\mathbb{R}} \frac{1}{2} \overline{f}(x) g(x) + \frac{1}{2} \overline{f}(x) g(-x) \, dx = \langle f, Pg \rangle$$

Example 8.9.

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Let  $\mathcal{H} = L^2(\mathbb{T})$ . Define  $Pf = \frac{1}{2\pi} \int_{\mathbb{T}} f \, dx, \, P : \mathcal{H} \to \mathcal{H}$ .

Given: 
$$f = \sum_{n \in \mathbb{Z}} \hat{f}(n) e^{inx}$$
  
Then:  $Pf = \hat{f}(0)$ 

• Idempotent:  $P^2 = P$  since Pf is a constant, and P1 = 1

• Self-Adjoint: 
$$\langle Pf,g\rangle = \int \overline{\left[\frac{1}{2\pi}\int f\,dx\right]}g\,dx = \frac{1}{2\pi}\int \overline{f}\,dx\int g\,dx = \langle f,1\rangle \cdot \frac{1}{2\pi}\int g\,dx = \langle f,Pg\rangle$$

ran P = constant functions =< 1 > (space spanned by 1) ker P = functions with zero mean (i.e.  $\hat{f}(0) = 0$ ) ran  $P \perp \ker P$ 

#### **Example 8.10.** Fourier Projections Notes 2/7/11

We can define the orthogonal projection of f onto the Nth partial sum of its Fourier series:

$$P_N f = \sum_{|n| \le N} \hat{f}(n) e^{inx}$$

Similarly, we can define the projection onto the positive n part of its Fourier series:

$$Pf = \sum_{n=0}^{\infty} \hat{f}(n)e^{inx}$$
$$(I-P)f = \sum_{n=-\infty}^{-1} \hat{f}(n)e^{inx}$$

### Example 8.11.

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Let  $\mathcal{H} = L^2(\mathbb{R})$ . If  $A \subset \mathbb{R}$  is some Lebesgue measurable set, define

$$\chi_A(x) = \begin{cases} 1 & x \in A \\ 0 & x \notin A \end{cases}$$

Then

$$P_A f = \chi_A f$$

is an orthogonal projection of  $L^2(\mathbb{R})$  onto the subspace of functions with support contained in  $\overline{A}$ .

#### 8.2 The Dual of a Hilbert Space

#### Theorem 8.12. Riesz Representation Theorem page 191 and Notes 2/7/11

Given: a Hilbert space  $\mathcal{H}$ , its dual space  $\mathcal{H}^* = \mathcal{B}(\mathcal{H}, \mathbb{C})$  (the set of bounded linear maps  $\varphi : \mathcal{H} \to \mathbb{C}$ with  $\|\varphi\|_{\mathcal{H}^*} = \sup \frac{|\varphi(x)|}{\|x\|} < \infty$ .

Every  $\varphi \in \mathcal{H}^*$  can be given by  $\varphi(x) = \langle y, x \rangle$  for some  $y \in \mathcal{H}$ , and  $\|\varphi\| = \|y\|$ . Conversely, every  $y \in \mathcal{H}$  corresponds to a  $\varphi \in \mathcal{H}^*$ . The map  $J: \varphi \mapsto y$  is an isometric, antilinear isomorphism of  $\mathcal{H}^*$ onto  $\mathcal{H}$ .

Antilinear: 
$$J(\varphi + \psi) = J(\varphi) + J(\psi)$$
  
 $J(\lambda \varphi) = \overline{\lambda} J(\varphi)$ 

Proof.

• Suppose  $\varphi \in \mathcal{H}^*$ . We want to find  $y \in \mathcal{H}$  such that  $\varphi(x) = \langle y, x \rangle$ 

- Suppose  $\varphi \neq 0$ . Then ker  $\varphi \neq \mathcal{H}$  and ker  $\varphi$  is closed because  $\varphi$  is bounded
- There exists  $z \in (\ker \varphi)^{\perp}$  (by the Projection Theorem)
- Consider  $P: \mathcal{H} \to \mathcal{H}, Px = \frac{\varphi(x)}{\varphi(z)} Pz$ . Claim: this is an orthogonal projection.
  - **Idempotent:**  $P^2 x = P\left(\frac{\varphi(x)}{\varphi(z)}z\right) = \frac{\varphi(x)}{\varphi(z)}Pz = \frac{\varphi(x)}{\varphi(z)}z$  (since Pz = z)
  - Self-Adjoint:  $\langle x, Py \rangle = \langle Px, y \rangle$
- $\mathcal{H} = \operatorname{ran} P \oplus \ker P$ ,  $\operatorname{ran} P = \langle z \rangle$ ,  $\ker P = \ker \varphi$
- $x \in \mathcal{H}$ ,  $x = \alpha z + w$ ,  $w \in \ker \varphi$ ,  $\alpha = \frac{\langle z, x \rangle}{\|z\|^2}$
- $\varphi(x) = \alpha \varphi(z) = \frac{\langle z, x \rangle}{\|z\|^2} \varphi(z) = \langle y, x \rangle, \quad y = \frac{\overline{\varphi}(z)}{\|z\|^2} z$

# 8.3 The Adjoint of an Operator

# **Definition 8.13.** *Adjoint* page 193 and Notes 2/9/11

Given a bounded linear map  $A \in \mathcal{B}(\mathcal{H})$ , its *adjoint*  $A^* \in \mathcal{B}(\mathcal{H})$  ( $\leftarrow$  proved in Proposition 8.15) is the linear map that satisfies

 $\langle x, Ay \rangle = \langle A^*x, y \rangle \quad \forall \ x, y \in \mathcal{H}$ 

# **Remark 8.14.** *Adjoint: Existence and Uniqueness* page 193 and Notes 2/9/11

To define  $A^*$  such that  $A^*x = z$ , consider  $\varphi_x : \mathcal{H} \to \mathbb{C}, \varphi_x(y) = \langle x, Ay \rangle$ . Then

$$\|\varphi_x(y)\| \le \|x\| \|Ay\| \le \|x\| \|A\| \|y\|$$
$$\|\varphi_x\| \le \|A\| \|x\|$$

So  $\varphi_x$  is a bounded linear functional. By the Riesz Representation Theorem, there is a unique  $z \in \mathcal{H}$  such that

$$\varphi_x(y) = \langle z, y \rangle$$

Define  $A^*x = z$ . Then

$$\langle x, Ay \rangle = \varphi_x(y) = \langle z, y \rangle = \langle A^*x, y \rangle \langle x, Ay \rangle = \langle A^*x, y \rangle \quad \forall \ x, y \in \mathcal{H}$$

**Proposition 8.15.** Notes 2/9/11

If  $A \in \mathcal{B}(\mathcal{H})$  then  $A^* \in \mathcal{B}(\mathcal{H})$  and

(1)  $||A^*|| = ||A||$ (2)  $||A||^2 = ||A^*A||$ 

(See also Corollary 8.34.)

Proof.

 $||A^*|| = \sup_{\|x\|=1} ||A^*x|| \quad \text{(See Lemma 8.26 in the book)}$  $= \sup_{\|x\|=\|y\|=1} |\langle y, A^*x \rangle| = \sup_{\|x\|=\|y\|=1} |\langle Ay, x \rangle| = \sup_{\|y\|=1} ||Ay|| = ||A||$ 

$$||A||^{2} = \sup_{\|x\|=1} ||Ax||^{2} = \sup_{\|x\|=1} |\langle Ax, Ax \rangle| = \sup_{\|x\|=1} |\langle x, A^{*}Ax \rangle|$$
  

$$\leq ||A^{*}A|| \quad (\text{See Corollary 8.27 in the book})$$
  

$$|A^{*}A|| \leq ||A^{*}|| ||A|| = ||A||^{2}$$
  

$$||A^{*}A|| = ||A||^{2}$$

Remark 8.16.

Notes 2/9/11

 $\mathcal{B}(\mathcal{H})$  is a  $C^*$ -algebra.

 $||AB|| \le ||A|| ||B|| \qquad *: \mathcal{B}(\mathcal{H}) \to \mathcal{B}(\mathcal{H}), \; ** = \text{identity} \qquad ||A^*|| = ||A||$ 

Remark 8.17. *Generalizations* Notes 2/9/11

- 1. Given:  $A : \mathcal{H} \to K, \ A^* : K \to \mathcal{H}$ , where  $\mathcal{H}, K$  are Hilbert spaces.  $\langle x, Ay \rangle_K = \langle A^*x, y \rangle_H \ \forall \ y \in \mathcal{H}, \ x \in K$  $A^*$  is the Hilbert space adjoint.
- 2. Given:  $A: X \to Y, A': Y' \to X'$ , where X, Y are Banach spaces and X' is the dual space of X.  $\langle \psi, Ax \rangle_{Y \times Y'} = \langle A'\psi, x \rangle_{X \times X'} \ \forall x \in X, \ \psi \in Y'$ 
  - A' is the dual operator or Banach space adjoint.

#### Example 8.18.

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Let  $\mathcal{H} = \mathbb{C}^n$ . Then  $A : \mathbb{C}^n \to \mathbb{C}^n$  is given by a matrix  $(a_{ij})$ .

$$y_i = \sum_{j=1}^n a_{ij} x_j, \qquad x = (x_1, \dots, x_n), \quad y = (y_1, \dots, y_n)$$
$$\langle x, y \rangle = \sum_{i=1}^n x_i y_i$$
$$\langle x, Ay \rangle = \sum_{i=1}^n \overline{x_i} \left( \sum_{j=1}^n a_{ij} y_j \right) = \sum_{j=1}^n \left( \overline{\sum_{i=1}^n \overline{a_{ij}} x_i} \right) y_j$$
$$= \langle A^* x, y \rangle$$

If 
$$z = A^* x$$
  
 $z_j = \sum_{i=1}^n \overline{a_{ij}} x_i = \sum_{j=1}^n \overline{a_{ji}} x_j$ 

- $A^*$  has matrix  $(\overline{a_{ji}})$ , which is the conjugate transpose of  $(a_{ij})$
- $(A^*A)$  is Hermitian, positive definite
- $(A^*A)^* = (A^*A)^* = A^*A$
- $\langle x, A^*Ax \rangle = \langle Ax, Ax \rangle \ge 0$
- $A^*A$  has orthogonal eigenvectors that form a basis of  $\mathbb{C}^n$  with eigenvalues  $\mu_1, \mu_2, \ldots, \mu_n \geq 0$
- $||A^*A|| = \max_{1 \le j \le n} |\mu_j| = \sigma(A^*A) =$ the spectral radius of  $A^*A$
- $||A|| = \sqrt{\sigma(A^*A)}$

#### Example 8.19.

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Let  $\mathcal{H} = L^2([0,1]), \ \langle f,g \rangle = \int_0^1 \overline{f(x)}g(x) \, dx.$ Define the integral operator  $K: L^2([0,1]) \to L^2([0,1])$  by

$$Kf(x) = \int_0^1 k(x, y) f(y) \, dy, \qquad k : [0, 1] \times [0, 1] \to \mathbb{C}$$

(Note: k(x, y) is the kernel of the integral operator K. It is not related to the null space.) **Ex:** Assume that k is *Hilbert-Schmidt*: k is measurable on  $[0, 1] \times [0, 1]$  and

$$||K||^{2} \leq \int_{0}^{1} \int_{0}^{1} |k(x,y)|^{2} \, dx \, dy < \infty$$

$$\begin{split} \langle f, Kg \rangle &= \int_0^1 \overline{f(x)} \left( \int_0^1 k(x, y) g(y) \, dy \right) \, dx \\ &= \int_0^1 \left( \overline{\int_0^1 f(x) \overline{k(x, y)} \, dx} \right) g(y) \, dy \\ &= \langle K^* f, g \rangle \end{split}$$

Since

$$K^*f(y) = \int_0^1 \overline{k(x,y)} f(y) \, dx$$
$$K^*f(x) = \int_0^1 \overline{k(y,x)} f(y) \, dy$$

Thus,  $K^*$  is an integral operator with conjugate transpose level of k.

#### Example 8.20.

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Recall the right and left shift operators, respectively:

 $S(x_1, x_2, x_3, \ldots) = (0, x_1, x_2, x_3, \ldots) \qquad T(x_1, x_2, x_3, \ldots) = (x_2, x_3, \ldots)$ 

T is the adjoint of S, i.e.  $T = S^*$ . Also,  $S = T^*$ .

**Example 8.21.** Solvability of Linear Equations Notes 2/11/11

Consider  $A : \mathcal{H} \to \mathcal{H}, Ax = y$ . Suppose for some  $y \in \mathcal{H}$  we have a solution for  $x \in \mathcal{H}$ .

Let  $z \in \ker A^*$ . Then

 $\langle z, Ax \rangle = \langle A^*z, x \rangle = \langle z, y \rangle$ 

Thus, a necessary condition for solvability is that  $y \perp z \forall z \in \ker A^*$ , i.e.  $y \perp \ker A^*$ .

Theorem 8.22. page 194 and Notes 2/11/11

If 
$$A \in \mathcal{B}(\mathcal{H})$$
, then  $\mathcal{H} = \overline{\operatorname{ran} A} \oplus (\ker A^*)$ , and

 $\overline{\operatorname{ran} A} = (\ker A^*)^{\perp} \qquad \ker A = (\operatorname{ran} A^*)^{\perp}$ 

*Proof.* From Example 8.21, if  $y \in \operatorname{ran} A$  then  $y \in (\ker A^*)^{\perp}$ .

 $\operatorname{ran} A \subset (\ker A^*)^{\perp}$  $\operatorname{ran} A \subset (\ker A^*)^{\perp} \qquad \text{since orthogonal complements are closed}$ 

If  $y \in (\operatorname{ran} A)^{\perp}$  then

$$\langle Ax, z \rangle = 0 \ \forall \ x \in \mathcal{H} \\ \langle x, A^*y \rangle = 0 \ \forall \ x \in \mathcal{H}$$

This implies that  $A^*y = 0$ , so  $y \in \ker A^*$ .

$$(\operatorname{ran} A)^{\perp} \subset \ker A^*$$
  
 $\overline{\operatorname{ran} A} = (\operatorname{ran} A)^{\perp \perp} \supset (\ker A^*)^{\perp}$ 

Corollary 8.23. page 195 and Notes 2/11/11

If  $A \in \mathcal{B}(\mathcal{H})$  has closed range (ran A is a closed linear subspace), then Ax = y is solvable iff  $y \perp \ker A^*$ .

# Example 8.24.

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If  $\mathcal{H}$  is finite dimensional, or A has finite rank, then ran A is closed and Corollary 8.23 applies.

# Example 8.25. page 196 and Notes 2/11/11

Recall the left (T) and right (S) shift operators.  $S^* = T, T^* = S$ .

1.  $\mathcal{H} = \overline{\operatorname{ran} S} \oplus \ker S^* = \overline{\operatorname{ran} S} \oplus \ker T$ 

- 2.  $\mathcal{H} = \overline{\operatorname{ran} T} \oplus \ker T^* = \overline{\operatorname{ran} T} \oplus \ker S$
- ran  $S = \{(x_1, x_2, \ldots) \in \ell^2 \mid x_1 = 0\}$  ran  $T = \ell^2(\mathbb{N})$
- $\ker S = \{0\}$

• ker  $T = \{(x_1, 0, 0, 0, \ldots) \mid x_1 \in \mathbb{C}\}$ 

Sx = y is solvable iff  $y \perp \ker T$ , and the solution is unique.

Tx = y is solvable for all  $y \in \ell^2(\mathbb{N})$ , but the solution is not unique.

#### Example 8.26.

Notes 2/11/11

#### 8.4 Self-Adjoint and Unitary Operators

**Definition 8.27.** Self-Adjoint page 197 and 2/14/11

A bounded operator  $A: \mathcal{H} \to \mathcal{H}$  on a Hilbert space  $\mathcal{H}$  is *self-adjoint* if  $A^* = A$ .

Equivalently, A is *self-adjoint* iff

 $\langle x, Ay \rangle = \langle Ax, y \rangle \qquad \forall \ x, y \in \mathcal{H}$ 

#### **Example 8.28.** Self-Adjoint Operators Notes 2/14/11

- 1.  $A : \mathbb{C}^n \to \mathbb{C}^n, \ [A]^* = [A]$  $A : \mathbb{R}^n \to \mathbb{R}^n, \ [A]^T = [A]$
- 2.  $\mathcal{H} = L^2(\mathbb{R})$ . Suppose  $a : \mathbb{R} \to \mathbb{C}$  is bounded and measurable. Define  $M : \mathcal{H} \to \mathcal{H}$ , Mf = af.  $\|Mf\|_2 \le \|a\|_{\infty} \|f\|_2$ .  $M^*f = \overline{a}f$ ,  $M^* = M$  if  $a : \mathbb{R} \to \mathbb{R}$ .
- 3. Orthogonal projections:  $P^2 = P = P^*$  (self-adjoint)
- 4. Given  $T \in \mathcal{B}(\mathcal{H})$ ,  $A = T^*T$  is self-adjoint. T = A + iB,  $A = \frac{1}{2}(T^* + T)$ ,  $B = \frac{1}{2i}(T^* - T)$  $A^* = A$ ,  $B^* = B$
- 5. The shift operators are NOT self-adjoint because  $S^* = T \neq S$

**Definition 8.29.** *Bilinear Forms, Sesquilinear* page 197 and Notes 2/14/11

Let  $A: \mathcal{H} \to \mathcal{H}$  be a bounded linear operator. We define the *bilinear form*  $a: \mathcal{H} \times \mathcal{H} \to \mathbb{C}$  by

$$a(x,y) = \langle x, Ay \rangle$$

We say that a is *sesquilinear* because

$$a(x, \lambda y + \mu z) = \lambda a(x, y) + \mu a(x, z)$$
$$a(\lambda x + \mu y, z) = \overline{\lambda} a(x, z) + \overline{\mu} a(x, z)$$

**Definition 8.30.** *Hermitian Symmetric & Symmetric* page 197 and Notes 2/14/11

Suppose A is self-adjoint. Then

$$\langle x, Ay \rangle = \langle Ax, y \rangle = \overline{\langle y, Ax \rangle}$$
  
 $a(x, y) = \overline{a(x, y)}$ 

We say that a is *Hermitian symmetric*. In the real case, we have a(x, y) = a(y, x), and we say that this is *symmetric*.

**Definition 8.31.** *Quadratic Form* page 197 and Notes 2/14/11

Given  $A: \mathcal{H} \to \mathcal{H}$ , we define the quadratic form  $q: \mathcal{H} \to \mathbb{C}$  by

 $q(x) = \langle x, Ax \rangle = a(x, x)$ 

If A is self-adjoint, then  $a(x, x) = \overline{a(x, x)}$ , so a(x, x) is real for all  $x \in \mathcal{H}$ .

#### Definition 8.32. Positive, Positive Definite page 198 and Notes 2/14/11

A self-adjoint operator A is positive or positive definite if  $\langle x, Ax \rangle = a(x, x) > 0$  for all  $x \in \mathcal{H}, x \neq 0$ .

#### Theorem 8.33.

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If A is self-adjoint then

$$|A|| = \sup_{x \neq 0} \frac{|\langle x, Ax \rangle|}{\|x\|^2} = \sup_{\|x\|=1} |\langle x, Ax \rangle|$$

Note: compare this to  $||A|| = \sup_{x \to a} |\langle Ax, Ax \rangle|^{1/2}$  (see part 2 of Proposition 8.15). ||x|| = 1

Proof.

 $|\langle x, Ax \rangle| \le ||x|| ||Ax|| \le ||A|| ||x||^2 \qquad \text{(Cauchy-Schwarz)}$ 

Let  $\alpha = \sup_{\|x\|\neq 0} \frac{|\langle x, Ax \rangle|}{\|x\|^2} \le \|A\|$ . Then  $|\langle x, Ax \rangle| \le \alpha \|x\|^2 \le \|A\| \|x\|^2$ . The parallelogram law states that

$$\langle x, Ay \rangle = \frac{1}{4} \left\{ \langle x+y, A(x+y) \rangle - \langle x-y, A(x-y) \rangle - i \langle x+iy, A(x+iy) \rangle + i \langle x-iy, A(x-iy) \rangle \right\}$$

In general,

$$||A|| = \sup_{||x|| = ||y|| = 1} |\langle x, Ay \rangle|$$

and this does not require self-adjoint. If A is self-adjoint, the first 2 terms in the parallelogram law expression are real and the last 2 are imaginary. We can multiply y by  $e^{i\theta}$  so that  $e^{i\theta} \langle x, Ay \rangle = \langle x, Az \rangle$  is real, where  $z = ye^{i\theta}$ . Then we have

$$e^{i\theta} \langle x, Ay \rangle = \langle x, Az \rangle$$

$$= \frac{1}{4} \{ \langle x + z, A(x + z) \rangle - \langle x - z, A(x - z) \rangle \}$$

$$| \langle x, Ay \rangle | \leq \frac{1}{4} | \langle x + z, A(x + z) \rangle | + \frac{1}{4} | \langle x - z, A(x - z) \rangle |$$

$$\leq \frac{\alpha}{4} \left( ||x + z||^2 + ||x - z||^2 \right)$$

$$\leq \frac{\alpha}{2} \left( ||x||^2 + ||z|| \right) \quad \text{(by the parallelogram rule (not law))}$$

$$||A|| \leq \sup_{||x||=||y||=1} | \langle x, Ay \rangle | \leq \frac{\alpha}{2} (||x||^2 + ||y||^2) \leq \frac{\alpha}{2} (1 + 1) = \alpha$$

Corollary 8.34. page 199

> If A is a bounded operator on a Hilbert space then  $||A^*A|| = ||A||^2$ . If A is self-adjoint, then  $||A^2|| = ||A||^2.$

The proof follows directly from Proposition 8.15.

**Definition 8.35.** Unitary Operators pages 199 & 200 and Notes 2/14/11

An operator  $U: \mathcal{H} \to \mathcal{H}$  is *unitary* if

 $U^*U = UU^* = I$ , i.e.  $U^* = U^{-1}$ 

Note that

$$\langle Ux, Uy \rangle = \langle U^*Ux, y \rangle = \langle x, y \rangle$$

so U preserves norms and inner products. Furthermore, if  $\{e_n \mid n \in \mathbb{N}\}$  is an orthonormal basis of  $\mathcal{H}$ , then so is  $\{Ue_n \mid n \in \mathbb{N}\}$ .

# Example 8.36.

Notes 2/14/11

1.  $U: \mathbb{C}^2 \to \mathbb{C}^2$  with matrix

$$[U] = \begin{pmatrix} a & b \\ -\overline{b} & \overline{a} \end{pmatrix}, \qquad |a|^2 + |b^2| = 1, \quad a, b \in \mathbb{C}$$

In the real case,  $a = \cos \theta$ ,  $b = \sin \theta$ , and U is rotation by  $\theta$ .

2. The right shift operator S on  $\ell^2(\mathbb{N})$  is not unitary because

$$S^* = T, \qquad S^*S = I, \qquad SS^* = P \neq I$$

3. If  $A^* = A$  then  $U = e^{iA}$  is unitary, where

$$e^{iA} = I + (iA) + \dots + \frac{1}{n!}(iA)^n + \dots$$
$$U^* = e^{-iA}$$
$$U^*U = I$$

Example 8.37. *Quantum Mechanics* Notes 2/14/11

In quantum mechanics we have the Hamiltonian operator H, with  $H^* = H$ . We also have  $U(t) = e^{itH}$ ,  $U : \mathcal{H} \to K$ ,  $U^* : K \to \mathcal{H}$ . U is unitary if  $U^*U = I_H$  and  $UU^* = I_K$ . We say that 2 Hilbert spaces are *isometric* if they are unitarily equivalent.

Example 8.38. page 201 and Notes 2/14/11

$$\mathcal{F}: L^2(\mathbb{T}) \to \ell^2(\mathbb{Z})$$
 is unitary  
 $\mathcal{F}f = \hat{f}, \quad \hat{f}(n) = \frac{1}{\sqrt{2\pi}} \int_{\mathbb{T}} f(x) e^{-inx} dx$ 

**Definition 8.39.** Normal Operators Notes 2/16/11

If  $T: \mathcal{H} \to \mathcal{H}$  is a bounded linear operator on a Hilbert space  $\mathcal{H}$ , then T is normal if

$$[T^*, T] \equiv T^*T - TT^* = 0$$
 i.e.  $T^*T = TT^*$ 

Self-adjoint and unitary operators are normal.

Example 8.40. Notes 2/16/11

Notes 2/16/11

- 1. Self-adjoint and unitary operators are normal
- 2. The shift operators on  $\ell^2(\mathbb{N})$  are not normal
- 3. Any multiplication operator is normal

$$\begin{split} M: L^2(\mathbb{R}) &\to L^2(\mathbb{R}) \\ (Mf)(x) &= m(x)f(x), \qquad m \in L^\infty(\mathbb{R}) \\ M^*f &= \overline{m}f \\ M^*Mf &= \overline{m}mf = m\overline{m}f = MM^*f \end{split}$$

Special cases

- (a) If m is real-valued then  $M = M^*$ , so M is self-adjoint. For
- (b) For M to be unitary, we must have  $m = e^{i\theta}$ .

#### 8.6 Weak Convergence in a Hilbert Space

**Definition 8.41.** *Weak Convergence* page 204 and Notes 2/16/11

A sequence  $(x_n)$  in a Hilbert space  $\mathcal{H}$  converges weakly to  $x \in \mathcal{H}$ , written  $x_n \rightharpoonup x$ , if

 $\langle x_n, y \rangle \to \langle x, y \rangle \quad \forall \ y \in \mathcal{H}$ 

Compare to Distributional Convergence (Definition 11.2):  $T_n \rightharpoonup T$  in  $\mathcal{D}'$  if  $\langle T_n, \varphi \rangle \rightarrow \langle T, \varphi \rangle$ .

We write strong (norm) convergence as  $x_n \to x$  if  $||x_n - x|| \to 0$ .

Remark 8.43. Weak vs. Strong Convergence Notes 2/16/11

If  $x_n \to x$ , then  $x_n \rightharpoonup x$  because

 $|\langle x_n, y \rangle - \langle x, y \rangle| \le ||x_n - x|| ||y||$  (Cauchy-Schwarz)

In a finite dimensional space, the converse is true, but this is not the case in infinite dimensional spaces.

Weak convergence = component-wise convergence

**Example 8.44.** page 204 and Notes 2/16/11

Let  $\mathcal{H}$  be a separable Hilbert space and let  $\{e_n \mid n \in \mathbb{N}\}$  be a separable orthonormal basis. Then  $e_n \rightharpoonup 0$  as  $n \rightarrow \infty$  because

 $\langle e_n, y \rangle = y_n \to 0$  as  $n \to \infty$  because  $\sum |y_n|^2 < \infty$ 

But  $(e_n)$  doesn't converge strongly because

$$\|e_n - e_m\| = \sqrt{2} \quad \forall \ n \neq m$$

and so the sequence is not Cauchy and hence not convergent.

# Example 8.45.

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Define an unbounded sequence  $(x_n)$  by  $x_n = ne_n$ . We know that

$$\langle x_n, e_m \rangle \to 0 \quad \Rightarrow \quad \langle x_n, y \rangle \to 0 \quad \text{as} \quad n \to \infty \quad \forall \ y = \sum_{m=1}^m c_m e_m$$

Let  $y_1 = \sum \frac{1}{m} e_m$ . Then

$$\langle x_n, y \rangle = \frac{1}{n} \cdot n = 1 \quad \forall \ n$$

Let  $y_2 = \sum \frac{1}{m^{3/4}} e_m \in \mathcal{H}$ . Then

$$\langle x_n, y \rangle = \frac{1}{n^{3/4}} \cdot n \to 0$$

Thus,  $(x_n)$  does not converge weakly.

# **Theorem 8.46.** Uniform Boundedness Theorem page 204

Suppose that  $\{\varphi_n : X \to \mathbb{C} \mid n \in \mathbb{N}\}$  is a set of functionals on a Banach space X such that the set of complex numbers  $\{\varphi_n(x) \mid n \in \mathbb{N}\}$  is bounded for each  $x \in X$ . Then  $\{\|\varphi_n\| \mid n \in \mathbb{N}\}$  is bounded.

### Theorem 8.47.

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If  $x_n \rightharpoonup x$  then  $\{ \|x_n\| \mid n \in \mathbb{N} \}$  is bounded.

*Proof.* Define  $\varphi_n : \mathcal{H} \to \mathbb{C}$  by  $\varphi_n(y) = \langle x_n, y \rangle$ . Then  $\varphi_n \in \mathcal{H}^*$ . By the uniform boundedness theorem (Theorem 8.46),

 $\begin{aligned} |\varphi_n(y)| &\leq M \quad \forall \ y \in \mathcal{H}, n \in \mathbb{N} \\ \{|\varphi_n(y)| \mid n \in \mathbb{N}\} \text{ is bounded for each } y \in \mathcal{H}, \text{ so } \{\|\varphi_n\| \mid n \in \mathbb{N}\} \text{ is bounded} \end{aligned}$ 

# Theorem 8.48.

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Let  $D \subset \mathcal{H}$  be a dense subset. Then  $x_n \to x$  iff (a)  $\{ \|x_n\| \mid n \in \mathbb{N} \}$  is bounded (b)  $\langle x_n, y \rangle \to \langle x, y \rangle \quad \forall y \in D$ 

#### Proposition 8.49. page 208 and Notes 2/16/11

If  $x_n \rightharpoonup x$ , then  $||x|| \le \liminf_{n \to \infty} ||x_n||$ 

Proof.

$$\|x\|^{2} = \langle x, x \rangle = \lim_{n \to \infty} \langle x_{n}, x \rangle \le \|x\| \liminf_{n \to \infty} \|x_{n}\|$$
$$\langle x_{n}, x \rangle \le \|x_{n}\| \|x\| \qquad (Cauchy-Schwarz)$$

Note: if  $a_n \leq b_n$ ,  $a_n \to a$ , then  $a \leq \liminf b_n$ .

$$||x_n - x||^2 = \langle x_n - x, x_n - x \rangle = ||x_n||^2 - \langle x, x_n \rangle - \langle x_n, x \rangle + ||x||^2$$

If  $x_n \rightharpoonup x$ , then  $||x_n|| \rightarrow ||x||$ , and

 $||x_n - x||^2 \to ||x||^2 - \langle x, x \rangle - \langle x, x \rangle + ||x||^2 = 0$ 

**Example 8.50.** *Example for Proposition 8.49* Notes 2/16/11

 $x_{1} = e_{1} \qquad x_{n} \rightharpoonup 0$   $x_{2} = 2e_{2} \qquad \|x_{n}\| = \begin{cases} 1 & n \text{ odd} \\ 2 & n \text{ even} \end{cases}$   $x_{4} = 2e_{4} \qquad \lim_{n \to \infty} \inf = 1$ ...

**Example 8.51.** Weak Convergence  $\Rightarrow$  Strong Convergence Notes 2/16/11

- (a) **Oscillation:** 
  - (1) Let  $\mathcal{H} = L^2(\mathbb{T}), f_n(x) = e^{inx} \to 0 \text{ as } n \to \infty$

*Proof.*  $||f_n|| = \sqrt{2\pi}$  is bounded, and  $\langle e^{inx}, \varphi \rangle \to 0$  as  $n \to \infty$  for all trig polynomials  $\varphi$ , and the trig polynomials are dense in  $L^2(\mathbb{T})$ .

- (2) Let  $\mathcal{H} = L^2(\mathbb{R})$ . Recall that  $C_C^{\infty}(\mathbb{R}) \subset L^2(\mathbb{R})$  are the smooth functions with compact support, and they are dense in  $L^2(\mathbb{R})$ . Then  $f_n \rightharpoonup f$  iff
  - i.  $||f|| \leq M$  (bounded)

ii.  $\int f_n \varphi \, dx \to \int f \varphi \, dx \, \forall \, \varphi \in C^\infty_C(\mathbb{R})$ 

Consider  $f_n(x) = \psi(x) \sin(n\pi x)$ , where  $\psi \in C_C^{\infty}(\mathbb{R}) \cap L^2(\mathbb{R})$ . Then  $f_n \to 0$  as  $n \to \infty$ , but  $f_n \neq 0$  as  $n \to \infty$ . (See proof below)

(b) Concentration: Consider

$$f_n(x) = \begin{cases} n^{1/2} & 0 < x < \frac{1}{n} \\ 0 & \text{otherwise} \end{cases}$$

i. 
$$||f_n||^2 = \int_0^{1/2} (n^{1/2})^2 dx = 1$$

ii.  $\forall \varphi \in C_C^{\infty}(\mathbb{R}), \quad \left| \int f_n \varphi \, dx \right| = \left| n^{1/2} \int_0^{1/n} \varphi \, dx \right| \le n^{1/2} \cdot \frac{1}{n} \|\varphi\|_{\infty} \to 0 \text{ as } n \to \infty$ So  $f_n \rightharpoonup 0$  as  $n \to \infty$ 

Does  $f_n$  converge strongly to 0? No, because  $||f_n|| = 1 \forall n$ . (See below for more details)

(c) Escape to Infinity:

$$f_n(x) = \begin{cases} 1 & n < x < n+1 \\ 0 & \text{otherwise} \end{cases}$$

i.  $||f_n||_{L^2} = 1$ , so  $f_n$  is bounded.

ii.  $\int f_n \varphi \, dx \to 0$  as  $n \to \infty \, \forall \varphi \in C_C^\infty(\mathbb{R})$ 

Thus,  $f_n \rightarrow 0$ , but  $f_n \not\rightarrow 0$  because  $||f_n|| = 1 \forall n$ .

Proof. (a2)

i. 
$$||f_n||^2 = \int \psi^2(x) \sin^2(n\pi x) \, dx \le \int \psi^2(x) \, dx \le ||\psi||^2$$

ii. Suppose  $\varphi \in C_C(\mathbb{R})$ .

$$\int f_n(x)\varphi(x) \, dx = \int \psi(x) \sin(n\pi x)\varphi(x) \, dx$$
$$= \int \frac{\cos(n\pi x)}{n\pi} \left[\varphi(x)\psi(x)\right]' \, dx$$
$$\left|\int f_n\varphi \, dx\right| \le \frac{1}{n\pi} \int (|\varphi\psi|)' \, dx$$
$$\le \frac{c}{n}$$

(IBP, no boundary terms because  $\varphi \in C_C(\mathbb{R})$ )

So  $\int f_n \varphi \, dx \to 0$  as  $n \to \infty$ , and thus  $f_n \rightharpoonup f$ .

Does  $(f_n)$  converge strongly? i.e., does  $f_n \to 0$ ? (see Remark 8.52) If  $\psi \neq 0$ , then

$$||f_n||^2 = \int \psi^2(x) \sin^2(n\pi x) \, dx = \int \psi^2(x) \cdot \frac{1}{2} \left[1 - \cos(2n\pi x)\right] \, dx \to \frac{1}{2} ||\psi||^2 \neq 0$$

In fact, if we set  $g_n = f_n^2 = [\psi(x)]^2 \sin^2(n\pi x)$ , then  $g_n \to \frac{1}{2}\psi^2(x)$  because

$$\int g_n(x)\varphi(x) \, dx = \int \psi^2(x) \sin^2(n\pi x)\varphi(x) \, dx$$
$$= \frac{1}{2} \int \psi^2 \varphi \, dx - \frac{1}{2} \int \varphi^2 \psi \cos(2\pi nx) \, dx$$
$$\to \frac{1}{2} \int \psi^2 \varphi \, dx$$

So  $g_n \rightharpoonup \frac{1}{2}\psi^2$ 

*Proof.* (b)

$$g_n = \begin{cases} n & 0 < x < \frac{1}{n} \\ 0 & \text{otherwise} \end{cases}$$

 $||g_n|| = \sqrt{n}, (g_n)$  is unbounded, so  $g_n \not\rightharpoonup g$ . In fact,  $g_n \rightharpoonup \delta \in \mathcal{D}'(\mathbb{R})$ .

$$h_n = \begin{cases} n^{1/4} & 0 < x < \frac{1}{n} \\ 0 & \text{otherwise} \end{cases}$$

 $||h_n|| = 0$ , and  $(h_n)$  is strongly and weakly convergent to 0.  $\frac{1}{2}$  is the critical value for  $L^2$ , and  $\frac{1}{p}$  is the critical value for  $L^p$ .

Remark 8.52.

If  $f_n \rightharpoonup f$  and  $f_n \rightarrow g$ , then we must have f = g because

Since  $\langle f,h\rangle = \langle g,h\rangle \ \forall h$ , we have that f = g.

#### 8.7 The Banach-Alaoglu Theorem

# Definition 8.53. Weakly Sequentially Compact

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A set  $K \subset \mathcal{H}$  is weakly sequentially compact if for any sequence  $(x_n) \subset K$  there exists a subsequence  $(x_{n_k})$  such that  $x_{n_k} \rightharpoonup x \in K$ .

Theorem 8.54. *Banach-Alaoglu Theorem* page 208 and Notes 2/23/11

Suppose that  $\mathcal{H}$  is a separable Hilbert space and  $\overline{B} = \{x \in \mathcal{H} \mid ||x|| \leq 1\}$  is the closed unit ball. Then  $\overline{B}$  is weakly sequentially compact.

#### Remarks

- 1.  $\overline{B}$  is not strongly compact if  $\mathcal{H}$  is infinite-dimensional. Ex:  $\{e_n\}$  is an orthonormal basis, but  $(e_n)$  has no convergent subsequence
- 2. This can be thought of as a replacement of the Heine-Borel theorem in the infinite-dimensional case

Proof. Let  $\{y_k \mid k \in \mathbb{N}\}$  be a dense subset of  $\mathcal{H}$ . Consider  $(\langle x_n, y_1 \rangle)_n \subset \mathbb{C}$ . By Cauchy-Schwarz,  $|\langle x_n, y \rangle \leq ||x_n|| ||y_1|| \leq ||y_1||$ , so the sequence is bounded, and thus there exists a subsequence of  $(x_n)$ , denoted  $(x_{n,1,k})_k = (x_{1,k})$  such that  $\langle x_{1,k}, y_1 \rangle$  converges as  $k \to \infty$ . Pick a subsequence  $(x_{2,k})$  of  $(x_{1,k})$  such that  $\langle x_{2,k}, y_2 \rangle$  converges as  $k \to \infty$ . Let  $x_j = x_{j,j}$  be the diagonal sequence. Then  $\langle x_j, y_n \rangle$  converges for every  $y_k$  as  $j \to \infty$  in this dense subset of  $\mathcal{H}$ . This defines a bounded linear functional F on  $D = \{y_k \mid k \in \mathbb{N}\}$ . By the Bounded Linear Transformation Theorem, this extends to a bounded linear functional  $\overline{F} : \mathcal{H} \to \mathbb{C}$  such that  $\overline{F}(y_k) = \lim_{j \to \infty} \langle x_j, y_k \rangle$  for all  $k \in \mathbb{N}$ . By the Riesz Representation Theorem, there exists  $x \in \mathcal{H}$  such that  $\langle x, y_k \rangle = \lim_{j \to \infty} \langle x_j, y_k \rangle$  for all  $k \in \mathbb{N}$ . Since  $\{y_k\}$  is dense in  $\mathcal{H}$  and  $||x|| \leq 1$ ,  $\langle x, y \rangle = \lim_{j \to \infty} \langle x_j, y_j \rangle$  for all  $y \in \mathcal{H}$ , and thus  $x_j \to x$ .  $||x|| \leq \liminf_{j \to \infty} ||x_j|| \leq 1$ , so  $x \in \overline{B}$ .

# Remark 8.55.

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- 1. We don't need  $\mathcal{H}$  to be separable (restrict to a closed subspace spanned by  $\{x_n\}$  which is separable)
- 2. Generalization to Banach spaces: the unit ball of  $X^*$  is weak-\* compact (equivalent to being weak compact if X is reflexive, i.e.  $X^{**} = X$ )

**Definition 8.56.** Weakly Sequentially Closed Notes 2/23/11

A set  $F \subset \mathcal{H}$  is weakly sequentially closed if whenever  $(x_n) \subset F$  is a sequence and  $x_n \rightharpoonup x$ , then  $x \in F$ .

**Example 8.57.** Weakly Closed  $\Rightarrow$  Strongly Closed Notes 2/23/11

Weakly closed implies strongly closed, but not conversely if  $\mathcal{H}$  is infinite-dimensional. For example, let

$$S = \{x \in \mathcal{H} \mid ||x|| = 1\}$$
$$\overline{B} = \{x \in \mathcal{H} \mid ||x|| \le 1\}$$

S is not weakly closed because  $(e_n) \subset S$ ,  $e_n \rightharpoonup 0 \notin S$ .  $\overline{B}$  is weakly closed because if  $x_n \rightharpoonup x$ , then  $||x|| \leq \liminf ||x_n||$ . The weak closure of S is  $\overline{B}$ .

**Definition 8.58.** Weakly Sequentially Lower Semicontinuous page 208 and Notes 2/23/11

A function  $f: D \subset \mathcal{H} \to \mathbb{R}$  is weakly sequentially lower semicontinuous if

$$x_n \rightarrow x \quad \Rightarrow \quad f(x) \leq \liminf_{n \rightarrow \infty} f(x_n)$$

Example:  $\|\cdot\| : \mathcal{H} \to \mathbb{R}$  is weakly sequentially lower semicontinuous.

**Remark 8.59.** Notes 2/23/11

Weakly sequentially lower semicontinuous implies strongly sequentially lower semicontinuous, but not conversely.

#### Theorem 8.60.

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Suppose that D is a weakly closed, bounded (in norm) subset in a Hilbert space  $\mathcal{H}$  and  $f: D \to \mathbb{R}$  is a weakly sequentially lower semicontinuous function. Then f is bounded from below  $(m = \inf_{x \in D} f(x) > -\infty)$  and there exists  $x \in D$  such that f(x) = m.

#### 8.8 Chapter Summary

We begin by defining what it means for a bounded linear operator P to be a projection (with "opposite" Q = I - P), and we explore relationship between projections and direct sum decompositions: P a projection  $\Leftrightarrow X = \operatorname{ran} P \oplus \ker P$ . We introduce orthogonal projections and show that they are bounded and self-adjoint. We explore the connection between orthogonal projections P ( $\Rightarrow \mathcal{H} = \operatorname{ran} P \oplus \ker P$ ) and direct sum decompositions ( $\mathcal{M}$  closed)  $\mathcal{H} = \mathcal{M} \oplus \mathcal{M}^{\perp}$  ( $\Rightarrow P$ ,  $\operatorname{ran} P = \mathcal{M}$ ,  $\ker P = \mathcal{M}^{\perp}$ ).

Recall from Chapter 5 that a linear functional is bounded iff it is continuous. We introduce the *Riesz Representation Theorem*: for all  $\varphi \in \mathcal{H}^*$ , there exists  $y \in \mathcal{H}$  such that  $\varphi(x) = \langle y, x \rangle$ . This gives us that all Hilbert spaces are self-dual:  $\mathcal{H}^{**} = \mathcal{H}$ . This is because the map  $J_1 : \mathcal{H} \to \mathcal{H}^*$  defined by  $J_1 y = \varphi_y$  identifies  $\mathcal{H}$  with its dual space,  $\mathcal{H}^*$ . Similarly, we can define a map  $J_2$  that identifies  $\mathcal{H}^*$  with its dual space,  $\mathcal{H}^{**}$ . Thus,  $\mathcal{H}$  and  $\mathcal{H}^{**}$  (and  $\mathcal{H}^{*}$ ) have the same cardinality. And since we know (Chapter 5) that for every  $x \in \mathcal{H}$  we can define a functional  $F_x \in \mathcal{H}^{**}$  by  $F_x(\varphi) = \varphi(x)$ , we therefore know that all linear functionals in  $\mathcal{H}^{**}$  are of this form.

We use the Riesz Representation Theorem to prove the existence of the *adjoint* of a bounded operator on a Hilbert space:  $\langle x, Ay \rangle = \langle A^*x, y \rangle$ . Examples:

• Matrix:  $A^* = A^T (\overline{A^T} \text{ if } A \text{ is complex})$ 

$$-\langle x, Ay \rangle = x^T Ay, \ \langle A^*x, y \rangle = (A^*x)^T y = x^T (A^*)^T y$$

- Integral operator  $Kf(x) = \int_0^1 k(x, y) f(y) \, dy$ :  $K^*f(x) = \int_0^1 \overline{k(y, x)} f(y) \, dy$
- Shift operators:  $S^* = T$ ,  $T^* = S$

We verify that for a bounded linear operator A, a solvability condition for Ax = y is that  $\langle y, z \rangle = 0$  for all  $z \in \ker A^* \Leftrightarrow \operatorname{ran} A \subset (\ker A^*)^{\perp}$ . We use this fact to prove that for a bounded linear operator A,

$$\overline{\operatorname{ran} A} = (\ker A^*)^{\perp}, \qquad \ker A = (\operatorname{ran} A^*)^{\perp}.$$

Equivalently,

$$\mathcal{H} = \underbrace{(\ker A^*)^{\perp}}_{\operatorname{ran} A} \oplus \underbrace{(\operatorname{ran} A)^{\perp}}_{\ker A^*}.$$

Next we have some definitions. We define what it means for a bounded linear operator to be *self-adjoint*, and we prove that for a bounded self-adjoint operator A,

$$||A|| = \sup_{||x||=1} |\langle x, Ax \rangle|, \qquad ||A^*A|| = ||A||^2.$$

Examples:

- A matrix is self-adjoint if it is symmetric (or Hermitian, if it is complex).
- An integral operator  $Kf(x) = \int_0^1 k(x, y) f(y) \, dy$  is self-adjoint if  $k(x, y) = \overline{k(y, x)}$

We say that an operator is *unitary/orthogonal* if it is invertible and  $\langle Ux, Uy \rangle_{\mathcal{H}_2} = \langle x, y \rangle_{\mathcal{H}_1} \Leftrightarrow U^*U = UU^* = I$ . We say that an operator is *normal* if  $T^*T = TT^*$ . (Self-adjoing and unitary operators are normal.)

Now we revisit weak convergence. For Hilbert spaces, the Riesz Representation Theorem gives us an equivalent definition:  $x_n \rightarrow x$  if  $\langle x_n, y \rangle \rightarrow \langle x, y \rangle \forall y \in \mathcal{H} \Leftrightarrow \varphi(x_n) \rightarrow \varphi(x) \forall \varphi \in \mathcal{H}^*$ . We mention 3 reasons why a sequence may converge weakly but not strongly: oscillation, concentration, and escape to infinity. We prove that for a weakly convergent sequence  $(x_n), ||x|| \leq \liminf ||x_n||$ . We also prove that if  $\lim ||x_n|| = ||x||$ , then  $(x_n)$  converges to x strongly. The Banach-Alaoglu Theorem tells us that the closed unit ball of a Hilbert space is weakly compact.

We define what it means for a function to be *convex*, and we say a few words about *lower semicontinuous* functions. We finish the chapter with *Mazur's Theorem*, which tells us that if  $x_n \rightharpoonup x$ , then there exists a sequence  $(y_n)$  of finite convex combinations of  $\{x_n\}$  that converges strongly to x.

# 9 The Spectrum of Bounded Linear Operators

#### 9.0 Introduction

#### Remark 9.1.

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Consider the following initial boundary value problem for a variable coefficient, linear equation:

$$u_t = u_{xx} - q(x)u \qquad 0 < x < 1, \ t > 0, u(0,t) = 0, \ u(1,t) = 0 \qquad t \ge 0, u(x,0) = f(x) \qquad 0 \le x \le 1$$

Using separation of variables, we assume

$$u(x,t) = \sum_{n=1}^{\infty} a_n(t)u_n(x)$$

where  $\{u_n \mid n \in \mathbb{N}\}$  is an orthonormal basis of  $L^2([0,1])$ . We find that

$$\frac{da_n}{dt} = -\lambda_n a_n n$$

and the  $u_n$  satisfy

$$-\frac{d^2u_n}{dx^2} + qu_n = \lambda_n u_n$$

Then the  $u_n$  are eigenvectors of the linear operator A. Thus,  $Au_n = \lambda_n u_n$ , where A is defined by

$$Au = -\frac{d^2u}{dx^2} + qu$$

We want a complete set of eigenvectors of A, or equivalently, to diagonalize A. This is an example of what we do in spectral theory.

#### 9.1 Diagonalization of Matrices

# Remark 9.2.

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The concept of the spectrum of an operator on a Banach/Hilbert space is a generalization of eigenvalues for matrices. Let  $A \in \mathcal{B}(X)$ . When dim  $X < \infty$  then we can identify it with a a matrix  $\tilde{A}$ . For any  $\lambda \in \mathbb{C}$  we have two possibilities:

- 1.  $\lambda I A$  is nonsingular  $\Leftrightarrow \det(\lambda I A) = 0 \Leftrightarrow (\lambda I A)^{-1}$  exists
- 2.  $\lambda I A$  is singular  $\Leftrightarrow$  there exists  $x_0$  such that  $(\lambda I A)x_0 = 0$ . Thus,  $Ax_0 = \lambda x_0$ ,  $\lambda$  is an eigenvalue, and  $x_0$  is an eigenvector.

What happens if dim  $X = \infty$ ???

#### 9.2 The Spectrum

Definition 9.3. Resolvent Set

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The resolvent set of a bounded operator A on a Banach space X is the set

 $\rho(A) = \{\lambda \in \mathbb{C} \mid (\lambda I - A) \text{ is invertible } \}$ (by the bounded inverse theorem)  $= \{\lambda \in \mathbb{C} \mid (\lambda I - A) \in \mathcal{B}(X) \}$  $= \{\lambda \in \mathbb{C} \mid (\lambda I - A) \text{ is } 1\text{-}1 \text{ and onto } \}$ 

**Definition 9.4.** *Spectrum* page 218 and Notes 3/2/11

The *spectrum* of A is the set

 $\sigma(A) = \mathbb{C} \setminus \rho(A)$ = { $\lambda \in \mathbb{C} \mid (\lambda I - A)$  is not invertible }

**Definition 9.5.** Point Spectrum, Continuous Spectrum, Residual Spectrum page 219 and Notes 3/2/11

In general,  $\sigma(A)$  can be expressed as  $\sigma(A) = \sigma_p(A) \cup \sigma_c(A) \cup \sigma_r(A)$ , where

- 1.  $\sigma_p(A) = \{\lambda \in \mathbb{C} \mid (\lambda I A) \text{ is not } 1\text{-}1 \}$  $\sigma_p(A)$  is called the *point spectrum* of A. In this case, since  $(\lambda I - A)$  is not 1-1, there exists  $x_0 \in \ker(\lambda I - A)$  such that  $(\lambda I - A)x_0 = 0 \Leftrightarrow Ax_0 = \lambda x_0$
- 2.  $\sigma_c(A) = \{\lambda \in \mathbb{C} \mid (\lambda I A) \text{ is 1-1 but not onto and } \overline{\operatorname{ran}(\lambda I A)} = X\}$  $\sigma_c(A)$  is called the *continuous spectrum* of A
- 3.  $\sigma_r(A) = \{\lambda \in \mathbb{C} \mid (\lambda I A) \text{ is 1-1 but not onto and } \overline{\operatorname{ran}(\lambda I A)} \neq X\}$  $\sigma_r(A)$  is called the *residual spectrum* of A

**Example 9.6.** Point, Continuous, and Residual Spectra Examples Notes 3/7/11

1. A matrix on  $\mathbb{C}^n$  has pure point spectrum

2.  $M: L^2([0,1]) \to L^2([0,1]), f \mapsto xf, \sigma(M) = [0,1]$  has pure continuous spectrum

3. Consider the right shift operator S on  $\ell^2(\mathbb{N})$ .  $\lambda = 0$  is in the residual spectrum

# Example 9.7.

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Consider the Banach space X = C([0,1]) with the  $\|\cdot\|_{\infty}$  norm. Define  $A : X \to X$  by Af(x) = xf(x). The boundedness of A follows exactly as in HW7 (even though  $X = L^2([0,1])$  on the HW, since we can take sup x = 1). Find  $\sigma(A)$ . Claim:  $\sigma(A) = \sigma_r(A) = [0,1]$ .

For any  $\lambda \in \mathbb{C}, f \in C([0,1])$ , we have

$$(\lambda I - A)f(x) = (\lambda - x)f(x) = 0$$

If  $\lambda \neq x$  then f(x) = 0. If  $\lambda \notin [0, 1]$  then  $\sigma_p = \emptyset$ .

For all  $\lambda \notin [0,1]$ , is  $(\lambda I - A)$  onto? For every  $g \in C([0,1])$ , we want f such that  $f(x)(\lambda - x) = g(x) \Rightarrow f(x) = \frac{g(x)}{\lambda - x} \in C([0,1])$ , since  $\lambda \notin [0,1]$  implies that  $\lambda - x \neq 0 \forall x \in [0,1]$ . Thus,  $(\lambda I - A)$  is onto, and we can conclude that  $\sigma(A) \subseteq [0,1]$ .

It will be enough to prove the claim to show that  $[0,1] \subseteq \sigma_r(A)$ . Why?  $[0,1] \subseteq \sigma_r(A) \subseteq \sigma(A) \subseteq [0,1]$ . Pick  $\lambda \in [0,1]$ . For every  $g \in \operatorname{ran}(\lambda I - A)$  we have that

$$g(x) = (\lambda - x)f(x) \text{ for some } f \in X = C([0, 1])$$
$$g(\lambda) = 0$$

So  $h(x) = 1 \notin \operatorname{ran}(\lambda I - A)$ , since  $g(\lambda) = 0 \neq 1$ . Therefore  $(\lambda I - A)$  is not onto.

If  $h \in \overline{\operatorname{ran}(\lambda I - A)}$  then there exists  $(g_n) \subset \operatorname{ran}(\lambda I - A)$  such that  $g_n \to h$ .  $h(\lambda = \lim_{n \to \infty} g_n(\lambda)(\lambda I - A) = 0$ . Thus,  $\mathbf{1} \notin \overline{\operatorname{ran}(\lambda I - A)}$ , so  $\lambda \in \sigma_r(A)$ .

### Example 9.8.

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Example 9.5 on page 219

**Definition 9.9.** *Resolvent* page 220 and Notes 3/4/11

For  $\lambda \in \rho(A)$ , we define the *resolvent* of A at  $\lambda$  to be

 $R_{\lambda} = (\lambda I - A)^{-1}, \qquad R_{\lambda} : \rho(A) \subset \mathbb{C} \to \mathcal{B}(\mathcal{H})$ 

# Example 9.10. Neumann Series

page 220 and Notes 3/4/11

If ||A|| < 1 then (I - A) is invertible and

$$(I - A)^{-1} = I + A + A^{2} + \dots$$

To show this, we define the partial sum:

$$S_N = I + A + A^2 + \ldots + A^N$$

Next, we show that the sequence of partial sums is Cauchy:

$$\begin{split} \|A^{M+1} + \ldots + A^N\| &\leq \|A^{M+1}\| + \ldots + \|A^N\| \leq \|A\|^{M+1} + \ldots + \|A\|^N \\ &\leq \sum_{n=M+1}^N \|A\|^n \end{split}$$

 $\sum_{n=1}^{\infty} < \infty$  if ||A|| < 1, so the partial sums are Cauchy. Thus,  $\sum_{n=0}^{\infty} A^n$  is Cauchy in  $\mathcal{B}(\mathcal{H})$ , and it converges since  $\mathcal{B}(\mathcal{H})$  is complete.

(See Remark 9.12.)

#### Example 9.11.

Notes 3/4/11

- 1. If  $|\lambda| > ||A||$  then  $\lambda \in \rho(A)$   $(\lambda I - A)^{-1} = \left[\lambda \left(I - \frac{A}{\lambda}\right)\right]^{-1} = \frac{1}{\lambda} \left(I - \frac{A}{\lambda}\right)^{-1}$  $\uparrow$  this exists if  $||A/\lambda|| < 1 \Rightarrow ||A|| < |\lambda|$
- 2. The resolvent set  $\rho(A)$  is open in  $\mathbb{C}$ Suppose  $\lambda_0 \in \rho(A)$ . We write:

$$\begin{aligned} (\lambda I - A) &= \lambda_0 I - A + (\lambda - \lambda_0) I = (\lambda_0 I - A) \left[ I + (\lambda - \lambda_0) (\lambda_0 I - A)^{-1} \right] \\ (\lambda I - A)^{-1} &= \left[ I + (\lambda - \lambda_0) (\lambda_0 I - A)^{-1} \right]^{-1} (\lambda_0 I - A)^{-1} \\ \uparrow \text{ exists if } |\lambda - \lambda_0| < \frac{1}{\|(\lambda_0 I - A)^{-1}\|} \end{aligned}$$

3. R<sub>λ</sub> : λ ↦ (λI - A)<sup>-1</sup> R<sub>λ</sub> is an operator-valued analytic function on the open set ρ(A) ⊂ C
4. σ(A) ≠ Ø

Remark 9.12.

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In Example 9.10, it is not necessary that ||A|| < 1 for  $(I - A)^{-1} = I + A + A^2 + \dots$  to converge. Rather, we require that  $\lim_{n \to \infty} ||A^n||^{1/n} < 1$ .

### Definition 9.13. Spectral Radius

page 220 and Notes 3/4/11

 $r(A) = \sup\{|\lambda| \mid \lambda \in \sigma(A)\}$  is the *spectral redius* of A. This is the radius of the smallest disc in  $\mathbb{C}$  centered at 0 that contains  $\sigma(A)$ . Also,  $r(A) \leq ||A||$ .

# Theorem 9.14.

page 220 and Notes 3/4/11

 $r(A) = \lim_{n \to \infty} \|A^n\|^{1/n}$  (and the limit exists)

*Proof.* Let  $a_n = \log ||A^n||$ . (If  $||A^n|| = 0$  for some n, i.e. A is nilpotent, then r(A) = 0.) Then

a

$$m+n = \log \|A^{m+n}\|$$
  

$$\leq \log \|A^n\| + \log \|A^n\|$$
  

$$\leq a_m + a_n \quad \text{(subadditive)}$$

We want to show that  $\lim_{n\to\infty} \frac{a_n}{n}$  exists, where  $\frac{a_n}{n} = \log ||A^n||^{1/n}$ . Fix n, m and write n = mp + q with  $0 \le q < m$ . Then we have

$$a_n = a_{mp+q} \le a_{mp} + a_q$$
$$\frac{a_n}{n} \le \frac{a_{mp}}{n} + \frac{a_q}{n}$$
Note that  $a_{mp} \le pa_m$ . Let  $n \to \infty$  with  $m$  fixed. Then  $\frac{p}{n} \to \frac{1}{m}$  as  $n \to \infty$ , and
$$\limsup \frac{a_n}{m} \le \frac{a_m}{m}$$
(9.1)

 $\limsup_{n \to \infty} \frac{n}{n} \le \frac{m}{m}$ 

Taking the limit of (9.1) as 
$$m \to \infty$$
, we obtain

$$\limsup_{n \to \infty} \frac{a_n}{n} \le \liminf_{m \to \infty} \frac{a_m}{m}$$

So  $\limsup_{n\to\infty} \frac{a_n}{n} = \limsup_{n\to\infty} \frac{a_n}{n}$ , and the sequence converges.

Example 9.15. Example for Theorem 9.14 Notes 3/4/11

$$A = \mu I \qquad \qquad \|A\| = |\mu| = r(A)$$
  
$$\lambda I - A = (\lambda - \mu) I \qquad \qquad \|A^n\|^{1/n} = |\mu|$$
  
$$\sigma(A) = \mu$$

**Corollary 9.16.** page 221 and Notes 3/4/11

If A is self-adjoint then r(A) = ||A||.

*Proof.*  $||A^2|| = ||A||^2$  and  $||A^{2^n}|| = ||A||^{2^n}$ , so  $\liminf_{n \to \infty} ||A^n||^{1/n} = ||A||$  by taking the subsequence  $n = 2^m$ .  $\Box$ 

#### 9.3 The Spectral Theorem for Compact, Self-Adjoint Operators

#### 9.3.1 Bounded, Self-Adjoint Operators

**Theorem 9.17.** page 222 and Notes 3/7/11

If A is bounded and self-adjoint, then every eigenvalue of A is real and eigenvectors with different eigenvalues are orthogonal.

Related to Theorem 9.21.

*Proof.* If  $Ax = \lambda x$ , then

$$\begin{aligned} \langle x, Ax \rangle &= \langle x, \lambda x \rangle = \lambda \|x\|^2 \\ \langle Ax, x \rangle &= \langle \lambda x, x \rangle = \overline{\lambda} \|x\|^2 \end{aligned}$$

If A is self-adjoint (and  $x \neq 0$ ), then  $\lambda = \overline{\lambda} \Rightarrow \lambda \in \mathbb{R}$ .

**Case:** A has pure point spectrum.

If  $Ax = \lambda x$  and  $Ay = \mu y$ ,  $x, y \neq 0$ ,  $\lambda \neq \mu$ , then

$$\begin{array}{l} \langle x, Ay \rangle = \mu \, \langle x, y \rangle \\ \langle Ax, y \rangle = \overline{\lambda} \, \langle x, y \rangle = \lambda \, \langle x, y \rangle \end{array} \right\} A = A^*, \text{ so } \mu \, \langle x, y \rangle = \lambda \, \langle x, y \rangle$$

If  $\lambda \neq \mu$ , then  $\langle x, y \rangle = 0$ , i.e.  $x \perp y$ .

What about the continuous and residual spectra?

$$\begin{split} \|(A - \lambda I)x\|^2 &= \langle (A - aI)x - ibx, (A - aI)x - ibx \rangle \qquad \text{where } \lambda = a + ib \\ &= \langle (A - aI)x, (A - aI)x \rangle + \underline{\langle -ibx, (A - aI)x \rangle} + \underline{\langle (A - aI)x, -ibx \rangle} + \langle -ibx, -ibx \rangle \\ &= \|(A - aI)x\|^2 + b^2 \|x\|^2 \\ &\geq b^2 \|x\|^2 \end{split}$$

Continuous Spectrum: See Proposition 9.18 and Remark 9.19. Residual Spectrum: See Proposition 9.20.

Proposition 9.18. page 223 and Notes 3/7/11

 $|\operatorname{Im} \lambda| \cdot ||x|| \le ||(A - aI)x||$ 

# Remark 9.19.

Notes 3/7/11

Proposition 9.18 says that if  $(A - \lambda I)x = y$ , then  $|\text{Im } \lambda| \cdot ||x|| \le ||y||$ . This means that if  $\lambda \in \mathbb{R}$ , we can estimate the solution, x, in terms of the RHS, y.

Applying this to the proof of Theorem 9.17, we see that if  $\lambda \in \mathbb{C} \setminus \mathbb{R}$ , it follows that

(a)  $(A - \lambda I)$  is 1-1 because if  $(A - \lambda I)x = 0$  then  $|\text{Im }\lambda|||x|| = 0 \Rightarrow x = 0$ .

(b)  $(A - \lambda I)$  has closed range. If  $y_n = (A - \lambda I)x_n, y_n \in ran(A - \lambda I), y_n \to y$ , then we can bound

$$\underbrace{\|x_m - x_n\|}_{\therefore \text{ Cauchy}} \le C \underbrace{\|y_m - y_n\|}_{\text{Cauchy}}$$

So  $x_n \to x$ ,  $(A - \lambda I)x = y$ , and  $y \in ran(A - \lambda I)$ . So if  $\lambda \in \mathbb{C} \setminus \mathbb{R}$ , then  $(A - \lambda I)$  is 1-1 with closed range, so there is no complex-valued continuous spectrum.

Proposition 9.20. page 224 and Notes 3/7/11

If A is bounded and self-adjoint, then the residual spectrum is empty.

*Proof.* If  $\lambda$  is in the residual spectrum, then there exists  $y \in \mathcal{H}$  such that  $\langle (A - \lambda I)x, y \rangle = 0 \quad \forall x \in \mathcal{H}$ , so  $y \perp \operatorname{ran}(A - \lambda I), y \neq 0$ . Since A is self-adjoint,  $\langle x, (A - \overline{\lambda}I)y \rangle = 0 \quad \forall x \in \mathcal{H}$ . This implies that  $(A - \overline{\lambda}I)y = 0$ , so y is an eigenvector of A with eigenvalue  $\overline{\lambda}$ . We have 2 cases:

- 1.  $\lambda \in \mathbb{C} \setminus \mathbb{R} \Rightarrow$  impossible (A has real eigenvalues)
- 2.  $\lambda \in \mathbb{R}$ . Then  $\lambda$  is in the point and residual spectra  $\Rightarrow$  impossible.

**Theorem 9.21.** page 223 and Notes 3/7/11

If A is a bounded, self-adjoint operator on a Hilbert space  $\mathcal{H}$ , then  $\sigma(A)$  is real and contained in the interval  $[-\|A\|, \|A\|]$ . The residual spectrum is empty.

Related to Theorem 9.17.

Proposition 9.22. page 223

If A is a bounded operator on a Hilbert space (not necessarily self-adjoint!) and  $\lambda \in \sigma_r(A)$ , then  $\overline{\lambda} \in \sigma_p(A^*)$ . In other words,  $\sigma_r(A) \subseteq \sigma_p(A^*)$ .

Bounded, self-adjoint operators have

- Spectral radius r(A) = ||A|| (See Corollary 9.16)
- Real eigenvalues (See Theorem 9.17)
- Orthogonal eigenvectors (See Theorem 9.17)
- Empty residual spectrum (See Proposition 9.20)

#### 9.3.2 Compact Operators

**Definition 9.24.** Compact Operator Notes 3/9/11

 $K: \mathcal{H} \to \mathcal{H}, D \in \mathcal{B}(\mathcal{H})$  is *compact* if it maps bounded sets to precompact sets.

Remark 9.25. Precompact Notes 3/9/11

Remember: a set is *precompact* if it is bounded and "almost" finite-dimensional.

Example 9.26. *The Hilbert Cube* page 230 and Notes 3/9/11

Let  $\mathcal{H} = \ell^2(\mathbb{N})$ . The Hilbert cube

$$C = \left\{ (x_1, x_2, \dots, x_n, \dots) \mid |x_n| \le \frac{1}{n} \right\}$$

is closed and precompact. Hence, C is a compact subset of  $\mathcal{H}$ .

**Example 9.27.** *Diagonal Operators Are Compact* page 230

The diagonal operator :  $\ell^2(\mathbb{N}) \to \ell^2(\mathbb{N})$  defined by

$$A(x_1, x_2, \dots, x_n, \dots) = (\lambda_1 x_1, \lambda_2 x_2, \dots, \lambda_n x_n, \dots)$$

is compact iff  $\lambda_n \to 0$  as  $n \to \infty$ .

#### **Example 9.28.** Compactness of Operators Notes 3/9/11

- 1. Any operator with finite rank (rank  $A = \dim \operatorname{ran} A$ ) is compact
- 2.  $I: \mathcal{H} \to \mathcal{H}$  is not compact if dim  $\mathcal{H} = \infty$
- 3.  $L^{2}([0,1]), Kf(x) = \int_{0}^{x} f(y) dy$  is a compact operator. If  $||f||_{L^{2}} \leq M$ , then

$$\left| \int_0^x f(y) \, dy \right| \le \int_0^1 |f(y)| \, dy \le \left( \int_0^1 |f(y)|^2 \, dy \right)^{1/2} \le M$$

Define  $F(x) = \int_0^x f(y) \, dy$ . Then

$$|F(x_2) - F(x_1)| = \left| \int_{x_1}^{x_2} f(y) \, dy \right| \le \left( \int_{x_1}^{x_2} 1 \cdot dy \right)^{1/2} \left( \int_{x_1}^{x_2} |f(y)|^2 \, dy \right)^{1/2} \le M |x_2 - x_1|^{1/2}$$

 $\{Kf \mid ||f|| \leq M\}$  is bounded and equicontinuous. Thus,  $H^2([0,1])$  is compactly embedded in  $L^2([0,1])$ . It follows that  $\{Kf \mid ||f||_{L^2} \leq M\}$  is precompact in C([0,1]) by Arzela-Ascoli, so it is precompact in  $L^2([0,1])$ .

If 
$$f(x) = \sum_{n=1}^{\infty} b_n \sin(n\pi x)$$
, then  $Kf(x) = \sum_{n=1}^{\infty} \frac{b_n}{n\pi} - \sum_{n=1}^{\infty} \frac{b_n}{n\pi} \cos(n\pi x)$ 

#### 9.3.3 **Compact**, Self-Adjoint Operators

#### **Remark 9.29.** page 223 and Notes 3/9/11

**Given:**  $A: \mathcal{H} \to \mathcal{H}, A$  is compact and self-adjoint,  $\mathcal{H}$  is a separable Hilbert space

We will prove:

- 1. A has at least one eigenvalue
- 2. If A leaves a subspace  $M \subset \mathcal{H}$  invariant  $(A: M \to M)$ , then A leaves  $M^{\perp}$  invariant, and  $\mathcal{H} = M \oplus M^{\perp}$

**Idea:** if we have  $A\varphi_n = \lambda_n \varphi_n$ , then we can get the largest eigenvalue by maximizing  $A(\sum c_n \varphi_n) =$  $\sum \lambda_n c_n \varphi_n$ .

Theorem 9.30.

page 225 and Notes 3/9/11

Suppose  $A: \mathcal{H} \to \mathcal{H}$  is compact and self-adjoint. Then A has an eigenvector with eigenvalue  $\lambda$  with  $\lambda = ||A||$  and/or  $\lambda = -||A||$ .

*Proof.* Recall: since A is self-adjoint,  $||A|| = \sup_{||x||=1} |\langle x, Ax \rangle|$ . Choose a sequence  $(x_n) \subset \mathcal{H}$  with  $||x_n|| = 1$ 

and  $\langle x_n, Ax_n \rangle \to \lambda$  as  $n \to \infty$ ,  $\lambda = \pm ||A||$ . Then we have

$$\begin{split} \|(A - \lambda I)x_n\|^2 &= \langle (A - \lambda I)x_n, (A - \lambda I)x_n \rangle \\ &= \langle Ax_n, Ax_n \rangle - 2\lambda \langle x_n, Ax_n \rangle + \lambda^2 \langle x_n, x_n \rangle \\ &= \underbrace{\|Ax_n\|^2}_{\leq \|A\|^2 \|x_n\|^2 = \lambda^2} -2\lambda \langle x_n, Ax_n \rangle + \lambda^2 \\ &\leq 2\lambda^2 - 2\lambda \langle x_n, Ax_n \rangle \qquad \to 0 \text{ as } n \to \infty \end{split}$$

So  $(A - \lambda I)x_n \to 0$  as  $n \to \infty$ , and thus  $x_n - \frac{1}{\lambda}Ax_n \to 0$  (assuming  $\lambda \neq 0$ , in which case ||A|| = 0 and everything is an eigenvalue). Since  $(x_n)$  is bounded  $(||x_n|| = 1 \forall n)$ ,  $Ax_n \to y$  by the compactness of A. So  $x_n \to \frac{y}{\lambda}$  and  $(A - \lambda I)y = 0$ .  $||y|| = \lambda \neq 0$ , since  $||x_n|| = 1$  and  $x_n \to y$ . So A has eigenvector y with eigenvalue  $\lambda$ .

# Proposition 9.31.

page 224 and Notes 3/9/11

- 1. Any nonzero eigenvalue of a compact, self-adjoint operator has a finite *multiplicity* (multiplicity  $\equiv$  the dimension of the eigenspace).
- 2. If  $\lambda_n$  is a sequence of eigenvalues and  $\lambda_n \to L$ , then we must have that L = 0.

**Theorem 9.32.** Spectral Theorem for Compact, Self-Adjoint Operators page 225 and Notes 3/11/11

If  $A : \mathcal{H} \to \mathcal{H}$  is a compact, self-adjoint operator on a Hilbert space  $\mathcal{H}$  then there is a finite or countably infinite sequence  $(\lambda_n)$  of nonzero real eigenvalues and orthogonal eigenvectors  $(\varphi_n)$  such that

$$A\varphi_n = \lambda_n \varphi_n$$

where  $|\lambda_1| \ge |\lambda_2| \ge \dots \ \lambda_n \to 0$  as  $n \to \infty$  if there are infinitely many  $\lambda_n$ 's and

$$Ax = \sum_{n} \lambda_n \langle \varphi_n, x \rangle \varphi_n$$
$$x = \sum_{n} \langle \varphi_n, x \rangle \varphi_n + n \quad \text{where } n \in \ker A, \quad \ker A \perp \underbrace{\langle \varphi_n \rangle}_{\text{span}}$$

Let  $P_n : \mathcal{H} \to \mathcal{H}$  be the orthogonal projection onto the eigenspace with eigenvalue  $\lambda_n$  (eigenvectors of bounded, self-adjoint operators are orthogonal; see Theorem 9.17). Then

$$A = \sum \lambda_n P_n$$

We are representing A as a sum of linear projections because  $\lambda_n \to 0$ , and so the sum converges uniformly.

*Proof.* To see that the sum converges uniformly to A, we compute

$$\|Ax - \sum_{n=1}^{N} \lambda_n P_n x\| = \sum_{n=N+1}^{\infty} |\lambda_n \langle \varphi_n, x \rangle \varphi_n|^2 \le |\lambda_{N+1}|^2 \|x\|^2$$

Also, if we let  $P_0$  be the orthogonal projection onto ker A, then

$$P_0 + \sum P_n = I$$

is strongly convergent. This is an example of what's called "resolution of the identity." Note that the  $\lambda_i$ 's gave us uniform convergence above. For bounded (and unbounded) self-adjoint operators with continuous spectrum we need to use resolutions of identity that involve integrals (instead of sums).

#### 9.4 Functions of Operators = Functional Calculus

# Definition 9.33. Function of an Operator page 232 and Notes 3/11/11If $f : \sigma(A) \subset \mathbb{C} \to \mathbb{C}$ is a bounded function, then we define $f(A) = \sum f(\lambda_n)P_n + f(0)P_0$ • f is uniformly convergent if $f(\lambda_n) \to 0$ as $n \to \infty$ • f is strongly convergent if $f(\lambda_n) \neq 0$ as $n \to \infty$ Note that $\sigma(A) = \{\lambda_n\} \cup \{0\}$ if dim $H = \infty$ • If there are finitely many $\lambda_n$ , then $0 \in \sigma_p(A)$

• If there are countably many  $\lambda_n$ , then  $0 \in \sigma_c(A)$ 

# Example 9.34.

Notes 3/11/11

Suppose A is a positive (see Definition 8.32), self-adjoint compact operator. Then

 $\langle x, Ax \rangle \ge 0$  implies  $\lambda_n \ge 0 \ \forall \ n$ 

We can define the positive square root of A as

$$\sqrt{A} = \sum \lambda_n^{1/2} P_n$$
$$\left(\sqrt{A}\right)^2 = \sum \lambda_n P_n = A$$

In general, if A is compact then

 $T = A^*A$  is positive and self-adjoint because  $\langle x, Tx \rangle = \langle x, A^*Ax \rangle = \langle Ax, Ax \rangle \ge 0$ 

$$\sqrt{T} = |A|, \qquad |A|^2 = T = A^*A$$

# **Definition 9.35.** *Polar Decomposition* page 217 and Notes 3/11/11

A = U|A|, where  $U : \operatorname{ran} |A| \to \operatorname{Im} A$  is a unitary operator

#### **Definition 9.36.** Fredholm Operator, Index Notes 3/11/11

A bounded operator  $A: \mathcal{H} \to \mathcal{H}$  is *Fredholm* if

- (a)  $\operatorname{ran} A$  is closed
- (b)  $\dim \ker A$  is finite
- (c)  $\operatorname{codim} \operatorname{ran} A$  is finite  $\Leftrightarrow \operatorname{dim} \ker A^*$  is finite

• codim ran  $A = \dim \ker A^*$  (recall that  $\mathcal{H} = \operatorname{ran} A \oplus \ker A^*$  when ran A is closed)

We define the *index* by

index  $A = \dim(\ker A) - \operatorname{codim}(\operatorname{ran} A) = \dim(\ker A) - \dim(\ker A^*)$ 

Example 9.37. Fredholm or not? Notes 3/11/11

- (a) I is Fredholm with index = 0
- (b)  $A(x_1, x_2, x_3, \ldots) \mapsto (x_1, \frac{1}{2}x_2, \frac{1}{3}x_3, \ldots)$  is not Fredholm because the range is not closed
- (c) The right shift operator, S, is Fredholm with index = -1

If A is Fredholm with index(A) = 0 then we have Fredholm alternative for solving the equation Ax = y, and there are 2 possibilities:

- 1. A is one-to-one and we can solve the equation for every  $y \in \mathcal{H}$
- 2. A is not one-to-one, and we can only solve the equation if  $y \perp \ker A^*$

Theorem 9.38. *Riesz-Schauder Theorem* Notes 3/11/11

If K is a compact, self-adjoint operator and  $\lambda \neq 0$  then  $A = \lambda I - K$  is Fredholm with index 0.

#### 9.5 Chapter Summary

$$U^{*}AU = U^{*}(AU) = U^{*} \left( A \begin{bmatrix} u_{1} & u_{2} & \cdots & u_{k} \end{bmatrix} \right)$$
  
=  $U^{*} \begin{bmatrix} Au_{1} & Au_{2} & \cdots & Au_{k} \end{bmatrix}$   
=  $U^{*} \begin{bmatrix} \lambda_{1}u_{1} & \lambda_{2}u_{2} & \cdots & \lambda_{k}u_{k} \end{bmatrix}$   
=  $\begin{bmatrix} \lambda_{1}e_{1} & \lambda_{2}e_{2} & \cdots & \lambda_{k}e_{k} \end{bmatrix}$  (because  $U^{*}u_{k} = e_{k}$ )  
=  $\begin{bmatrix} \lambda_{1} & 0 & \cdots & 0 \\ 0 & \lambda_{2} & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & \lambda_{k} \end{bmatrix} = D$ 

Operator	Spectrum	Point	Continuous	Residual
Bounded, Linear	Closed & Nonempty,			$\lambda \in \sigma_r(A) \; \Rightarrow \;$
	$r(A) = \lim \ A^n\ ^{1/n}$			$\overline{\lambda} \in \sigma_p(A^*)$
Bounded, Self-Adjoint	$\sigma(A) \subset [-\ A\ , \ A\ ]$	real	real	empty
	$r(A) = \ A\ $			
Compact, Self-Adjoint		$-\ A\  \in \sigma_p(A)$ or	$\sigma_c(A) = \{0\}$ or	empty
		$  A   \in \sigma_p(A)$	$\sigma_c(A) = \emptyset$	